# EXTENSION, SEPARATION AND ISOMORPHIC REVERSE ISOPERIMETRY 

ASSAF NAOR<br>Dedicated with awe to the memory of Jean Bourgain


#### Abstract

The Lipschitz extension modulus e $(M)$ of a metric space $M$ is the infimum over those $L \in[1, \infty]$ such that for any Banach space $\mathbf{Z}$ and any $\mathcal{C} \subseteq m$, any l-Lipschitz function $f: \mathcal{C} \rightarrow \mathbf{Z}$ can be extended to an $L$-Lipschitz function $F: M \rightarrow \mathbf{Z}$. Johnson, Lindenstrauss and Schechtman proved (1986) that if $\mathbf{X}$ is an $n$-dimensional normed space, then $\mathrm{e}(\mathbf{X}) \lesssim n$. In the reverse direction, we prove that every $n$-dimensional normed space $\mathbf{X}$ satisfies $\mathrm{e}(\mathbf{X}) \gtrsim n^{c}$, where $c>0$ is a universal constant. Our core technical contribution is a geometric structural result on stochastic clustering of finite dimensional normed spaces which implies upper bounds on their Lipschitz extension moduli using an extension method of Lee and the author (2005). The separation modulus of a metric space $(m, d m)$ is the infimum over those $\sigma \in(0, \infty]$ such that for any $\Delta>0$ there is a distribution over random partitions of $m$ into clusters of diameter at most $\Delta$ such that for every two points $x, y \in M$ the probability that they belong to different clusters is at most $\sigma d_{m}(x, y) / \Delta$. We obtain upper and lower bounds on the separation moduli of finite dimensional normed spaces that relate them to well-studied volumetric invariants (volume ratios and projection bodies). Using these connections, we determine the asymptotic growth rate of the separation moduli of various normed spaces. If $\mathbf{X}$ is an $n$ dimensional normed space with enough symmetries, then our bounds imply that its separation modulus is equal to $\operatorname{vr}\left(\mathbf{X}^{*}\right) \sqrt{n}$ up to factors of lower order, where $\operatorname{vr}\left(\mathbf{X}^{*}\right)$ is the volume ratio of the unit ball of the dual of $\mathbf{X}$. We formulate a conjecture on isomorphic reverse isoperimetric properties of symmetric convex bodies (akin to Ball's reverse isoperimetric theorem (1991), but permitting a non-isometric perturbation in addition to the choice of position) that can be used with our volumetric bounds on the separation modulus to obtain many more exact asymptotic evaluations of the separation moduli of normed spaces. Our estimates on the separation modulus imply asymptotically improved upper bounds on the Lipschitz extension moduli of various classical spaces. In particular, we deduce an improved upper bound on e $\left(\ell_{p}^{n}\right)$ when $p>2$ that resolves a conjecture of Brudnyi and Brudnyi (2005), and we prove that $\mathrm{e}\left(\ell_{\infty}^{n}\right)=\sqrt{n}$, which is the first time that the growth rate of $\mathrm{e}(\mathbf{X})$ has been evaluated $(\operatorname{as} \operatorname{dim}(\mathbf{X}) \rightarrow \infty)$ for any finite dimensional normed space $\mathbf{X}$.


Date: Last modified March 17, 2022.
Key words and phrases. Lipschitz extension, randomized clustering, convex geometry, local theory of Banach spaces, projection bodies, volume ratios, Wasserstein spaces, spectral geometry, Dirichlet eigenvalues, Cheeger sets, reverse isoperimetry.

Supported by NSF grant DMS-2054875, BSF grants 2010021 and 2018223, the Packard Foundation and the Simons Foundation. Part of this work was conducted under the auspices of the Simons Algorithms and Geometry (A\&G) Think Tank. An extended abstract Nao17a titled "Probabilistic clustering of high dimensional norms" that announces discrete and algorithmic aspects of parts of this work appeared in the proceedings of the 28th annual ACM-SIAM Symposium on Discrete Algorithms.

## Contents

1. Introduction ..... 3
1.1. Brief highlights of main results ..... 3
1.1.1. A conjectural isomorphic reverse isoperimetric phenomenon ..... 5
1.2. Basic notation ..... 9
1.3. Lipschitz extension ..... 9
1.4. A volumetric upper bound on the Lipschitz extension modulus ..... 13
1.5. A dimension-independent extension theorem ..... 18
1.6. Isomorphic reverse isoperimetry ..... 19
1.6.1. A spectral interpretation, reverse Faber-Krahn and the Cheeger space of a normed space ..... 21
1.6.2. Symmetries and positions ..... 24
1.6.3. Intersection with a Euclidean ball ..... 32
1.7. Randomized clustering ..... 33
1.7.1. Basic definitions related to random partitions ..... 34
1.7.2. Iterative ball partitioning ..... 34
1.7.3. Separation and padding ..... 36
1.7.4. From separation to Lipschitz extension ..... 37
1.7.5. Bounds on the separation and padding moduli of normed spaces ..... 37
1.7.6. Dimension reduction ..... 43
1.8. Consequences in the linear theory ..... 44
2. Lower bounds ..... 47
2.1. Proof of Theorem 13 ..... 47
2.2. Proof of (106) ..... 53
2.3. Hölder extension ..... 54
2.4. Justification of (25) ..... 57
2.5. Proof of the lower bound on $\operatorname{SEP}(\mathbf{X})$ in Theorem 3 ..... 58
2.6. Proof of the lower bound on $\mathrm{PAD}_{\delta}(\mathbf{X})$ in Theorem 68 ..... 61
2.7. Proof of Proposition 86 ..... 63
3. Preliminaries on random partitions ..... 64
3.1. Standard set-valued mappings ..... 66
3.2. Proximal selectors ..... 67
3.3. Measurability of iterative ball partitioning ..... 70
4. Upper bounds on random partitions ..... 71
4.1. Proof of Theorem 744 and the upper bound on $\mathrm{PAD}_{\delta}(\mathbf{X})$ in Theorem 68 ..... 72
4.1.1. Proof of the first inequality in (260) ..... 79
4.2. Proof of Theorem 80$]$ ..... 81
5. Barycentric-valued Lipschitz extension ..... 84
5.1. Notational preliminaries ..... 84
5.2. Refined extension moduli ..... 86
5.3. Barycentric targets ..... 89
5.4. Gentle partitions of unity ..... 90
5.5. The multi-scale construction ..... 92
6. Volume computations ..... 96
6.1. Direct sums ..... 96
6.2. Negatively correlated normed spaces ..... 118
6.3. Volume ratio computations ..... 120
7. Logarithmic weak isomorphic isoperimetry in minimum dual mean width position ..... 128
References ..... 131

## 1. Introduction

Our core technical contribution is a geometric structural result (stochastic clustering) for subsets of finite dimensional normed spaces. It provides new links between nonlinear questions in metric geometry and volumetric issues in convex geometry. An unexpected aspect of our statement is that it contradicts an impossibility result of the well-known work [CCG ${ }^{+} 98$ ] by Charikar, Chekuri, Goel, Guha and Plotkin in the computer science literature, thus leading to bounds that were previously thought to be impossible. This is reconciled in Section 1.7, where we explain the source of the error in [CG ${ }^{+}$98].

The aforementioned link opens up a vista that allows one to apply the extensive literature on the linear theory to important and well-studied nonlinear questions. It also raises new fundamental issues within the linear theory that we will only begin to address here. So, in order to fully explain both the history and the ideas and their consequences, we will start with a quick overview of some of our main results that assumes familiarity with standard concepts in the respective areas. We will then present a gradual and complete introduction to our work that specifies all of the necessary background.
1.1. Brief highlights of main results. Associate to every separable complete metric space $(m, d m)$ two bi-Lipschitz invariants e $(m), \operatorname{SEP}(T) \in(0, \infty]$ called, respectively, the Lipschitz extension modulus of $m$ and the separation modulus of $m$, that are defined as follows. The Lipschitz extension modulus of $m$ is the infimum over those $L \in(0, \infty]$ such that for every Banach space $\mathbf{Z}$ and every subset $\mathcal{C} \subseteq m$, every 1-Lipschitz function $f: \mathcal{C} \rightarrow \mathbf{Z}$ can be extended to a $\mathbf{Z}$-valued $L$-Lipschitz function that is defined on all of $m$. The separation modulus of $m$ is the infimum over those $\sigma \in(0, \infty]$ such that for any $\Delta>0$ there is a distribution over random partitions $\rrbracket^{1}$ of $m$ into clusters of diameter at most $\Delta$ such that for every two points $x, y \in M$ the probability that they belong to different clusters is at most $\sigma d_{m}(x, y) / \Delta$.

The question of estimating the Lipschitz extension modulus received great scrutiny over the past century; see Section 1.3 for an indication of (a small part of) the extensive knowledge on this topic. The separation modulus was introduced by Bartal in the mid-1990s and received a lot of attention in the computer science literature due to its algorithmic applications; see Section 1.7 .3 for the history. Its connection to Lipschitz extension was found by Lee and the author [LN04a, LN05], who proved that e $(m) \lesssim \mathrm{SEP}(m)$.

By a well-known theorem of Johnson, Lindenstrauss and Schechtman [JLS86], every normed space $\mathbf{X}$ satisfies $\mathrm{e}(\mathbf{X})=O(\operatorname{dim}(\mathbf{X}))$. Here we obtain a power-type lower bound on $\mathrm{e}(\mathbf{X})$ in terms of $\operatorname{dim}(\mathbf{X})$.
Theorem 1. There is a universal constant $c>0$ such that $(\mathbf{X}) \geqslant \operatorname{dim}(\mathbf{X})^{c}$ for every normed space $\mathbf{X}$.
Theorem 1 improves over the previously best-available bound $\mathrm{e}(\mathbf{X}) \geqslant e^{c \sqrt{\log \operatorname{dim}(\mathbf{X})}}$; see Remark 97 for the history of this question. Despite substantial efforts, the asymptotic growth rate (as $\operatorname{dim}(\mathbf{X}) \rightarrow \infty$ ) of $e(\mathbf{X})$ was not previously known (even up to lower order factors) for any sequence of normed spaces.
Theorem 2. For every $n \in \mathbb{N}$ we have ${ }^{2} \mathrm{e}\left(\ell_{\infty}^{n}\right)=\sqrt{n}$.
The previously best-known upper bound on $\mathrm{e}\left(\ell_{\infty}^{n}\right)$ was nothing better than the aforementioned general $O(n)$ bound of [JLS86]. Theorem 2 is just one instance of our asymptotically improved upper bounds on the Lipschitz extension moduli of many normed spaces of interest; we get e.g. the best-known bound when $\mathbf{X}=\ell_{p}^{n}$ for any $p>2$. Nevertheless, currently $\ell_{\infty}^{n}$ is essentially $y^{3}$ the only normed space whose Lipschitz extension modulus is known up to lower order factors (by Theorem (2), and the same question even for the Euclidean space $\ell_{2}^{n}$ remains a longstanding open problem; see Section 1.3 for more on this.

[^0]All of the upper bounds on the Lipschitz extension modulus that we obtain herein use the upper bound on the separation modulus that appears in Theorem 3 below. This theorem also contains a new lower bound on the separation modulus, which we will see shows that in several cases of interest our results are a sharp evaluation of the asymptotic growth rate of the separation modulus $\|^{4}$

Theorem 3. Let $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ and $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ be normed spaces whose unit balls satisfy $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$. Then

$$
\begin{equation*}
\operatorname{vr}\left(\mathbf{X}^{*}\right) \sqrt{n} \lesssim \operatorname{SEP}(\mathbf{X}) \lesssim \frac{\operatorname{diam}_{\mathbf{X}^{*}}\left(\Pi B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} . \tag{1}
\end{equation*}
$$

In the left hand side of (1), $\operatorname{vr}\left(\mathbf{X}^{*}\right)$ is the volume ratio Sza78, STJ80] of the dual $\mathbf{X}^{*}$, i.e., it is the $n$ 'th root of the ratio of the volume of $B_{\mathbf{X}^{*}}$ and maximal volume of an ellipsoid that is contained in $B_{\mathbf{X}^{*}}$. In the right hand side of (1], $\Pi B_{\mathbf{Y}}$ is the projection body [Pet67] of $B_{\mathbf{Y}}$, and $\operatorname{diam}_{\mathbf{X}^{*}}(\cdot)$ denotes diameter with respect to the metric on $\mathbb{R}^{n}$ that is induced by $\mathbf{X}^{*}$. We will recall the definition of a projection body latel ${ }^{5}$ and it suffices to mention now that the mapping $K \mapsto \Pi K$, which is of central importance in convex geometry (see [BL88, Lut93, Gar06, Sch14] for an indication of the extensive literature on this topic), associates to every convex body $K \subseteq \mathbb{R}^{n}$ a convex body $\Pi K \subseteq \mathbb{R}^{n}$ that encodes isoperimetric properties of $K$.

A key contribution of Theorem3 is the role of the auxiliary normed space $\mathbf{Y}$, which appears despite the fact that we are interested in the separation modulus of $\mathbf{X}$. By substituting $\mathbf{Y}=\mathbf{X}$ into the right hand side of (1) one does get a meaningful estimate, and in particular the resulting bound is $O(n$ ), i.e., (1) implies the bound of [JLS86]. However, we will see that by introducing a suitable perturbation $\mathbf{Y}$ of $\mathbf{X}$, the second inequality in (1) can sometimes be significantly stronger than the special case $\mathbf{Y}=\mathbf{X}$. We will exploit this powerful degree of freedom heavily; its geometric significance is discussed in Section 1.4 .

The previously best-known upper and lower estimates on the separation moduli of normed spaces are due to [CCG ${ }^{+} 98$, where it was proved that $\operatorname{SEP}\left(\ell_{1}^{n}\right)=n$ and $\operatorname{SEP}\left(\ell_{2}^{n}\right)=\sqrt{n}$. By bi-Lipschitz invariance, this implies that any $n$-dimensional normed space $\mathbf{X}$ satisfies

$$
\begin{equation*}
\frac{n}{d_{\mathrm{BM}}\left(\ell_{1}^{n}, \mathbf{X}\right)} \lesssim \mathrm{SEP}(\mathbf{X}) \lesssim d_{\mathrm{BM}}\left(\ell_{2}^{n}, \mathbf{X}\right) \sqrt{n}, \tag{2}
\end{equation*}
$$

where $d_{\mathrm{BM}}(\cdot, \cdot)$ denotes the Banach-Mazur distance. Both of the bounds in (2) can be inferior to those that follow from Theorem 33. For example, suppose that $n=m^{2}$ for some $m \in \mathbb{N}$ and consider $\mathbf{X}=\ell_{\infty}^{m}\left(\ell_{1}^{m}\right)$. Then, $d_{\mathrm{BM}}\left(\mathbf{X}, \ell_{1}^{n}\right)=d_{\mathrm{BM}}\left(\mathbf{X}, \ell_{2}^{n}\right)=\sqrt{n}$ by the work (KS89] of Kwapien and Schütt. Therefore in this case (2) becomes $\sqrt{n} \lesssim \operatorname{SEP}(\mathbf{X}) \lesssim n$, while we will see that (1) implies that $\operatorname{SEP}(\mathbf{X})=n^{3 / 4}$.

The following corollary collects examples of applications of Theorem 3 that we will deduce herein.
Corollary 4 (examples of consequences of Theorem(3). The following statements hold for any $n \in \mathbb{N}$.

- For any $p \geqslant 1$, the separation modulus of $\ell_{p}^{n}$ satisfies

$$
\begin{equation*}
\operatorname{SEP}\left(\ell_{p}^{n}\right)=n^{\max \left\{\frac{1}{2}, \frac{1}{p}\right\}} . \tag{3}
\end{equation*}
$$

More generally, let $\left(\mathbf{E},\|\cdot\|_{\mathbf{E}}\right)$ be any n-dimensional normed space with a 1 -symmetric basis $e_{1}, \ldots, e_{n}$. Then, $\operatorname{SEP}(\mathbf{E})$ is equal to the following quantity up to lower order factors:

$$
\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right)
$$

[^1]- For any $p \geqslant 1$, the separation modulus of the Schatten von-Neumann trace class $\mathrm{S}_{p}^{n}$ on $\mathrm{M}_{n}(\mathbb{R})$ is

$$
\begin{equation*}
\operatorname{SEP}\left(\mathrm{S}_{p}^{n}\right)=n^{\max \left\{1, \frac{1}{2}+\frac{1}{p}\right\}+o(1)}=\operatorname{dim}\left(\mathrm{S}_{p}^{n}\right)^{\max \left\{\frac{1}{2}, \frac{1}{4}+\frac{1}{2 p}\right\}+o(1)} \tag{4}
\end{equation*}
$$

More generally, let $\left(\mathbf{E},\|\cdot\|_{\mathbf{E}}\right)$ be any $n$-dimensional normed space with a 1 -symmetric basis $e_{1}, \ldots, e_{n}$ and denote its unitary ideal by $\mathrm{S}_{\mathrm{E}}=\left(\mathrm{M}_{n}(\mathbb{R}),\|\cdot\|_{\mathrm{S}_{\mathrm{E}}}\right)$. Then, $\operatorname{SEP}\left(\mathrm{S}_{\mathrm{E}}\right)$ is equal to the following quantity up to lower order factors:

$$
\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) \sqrt{n}
$$

- For any $p, q \geqslant 1$, the separation modulus of the $\ell_{p}^{n}\left(\ell_{q}^{n}\right)$ norm on $\mathrm{M}_{n}(\mathbb{R})$ is

$$
\begin{equation*}
\operatorname{SEP}\left(\ell_{p}^{n}\left(\ell_{q}^{n}\right)\right)=n^{\max \left\{1, \frac{1}{p}+\frac{1}{q}, \frac{1}{2}+\frac{1}{p}, \frac{1}{2}+\frac{1}{q}\right\}}=\operatorname{dim}\left(\ell_{p}^{n}\left(\ell_{q}^{n}\right)\right)^{\max \left\{\frac{1}{2}, \frac{1}{2 p}+\frac{1}{2 q}, \frac{1}{4}+\frac{1}{2 p}, \frac{1}{4}+\frac{1}{2 q}\right\}} . \tag{5}
\end{equation*}
$$

- For any $p, q \geqslant 1$, the separation modulus of $\mathrm{M}_{n}(\mathbb{R})$ equipped with the operator norm $\|\cdot\|_{\ell_{p}^{n} \rightarrow \ell_{q}^{n}}$ from $\ell_{p}^{n}$ to $\ell_{q}^{n}$ is equal to the following quantity up to lower order factors:

$$
\left\{\begin{array}{lll}
n^{\frac{3}{2}-\frac{1}{\min [p, q]}} & \text { if } & p, q \geqslant 2, \\
n^{\frac{1}{2}+\frac{\max }{\max , q]}} & \text { if } & p, q \leqslant 2, \\
n^{\max \left\{1, \frac{1}{q}-\frac{1}{p}+\frac{1}{2}\right\}} & \text { if } & \text { if } \\
n^{\operatorname{mas}}, 2 \leqslant p,
\end{array}\right.
$$

- For any $p, q \geqslant 1$, the separation modulus of the projective tensor product $\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}$, i.e., the norm on $\mathrm{M}_{n}(\mathbb{R})$ whose unit ball is the convex hull of $\left\{\left(x_{i} y_{j}\right) \in \mathrm{M}_{n}(\mathbb{R}) ;\left(x_{1}, \ldots, x_{n}\right) \in B_{\ell_{p}^{n}} \wedge\left(y_{1}, \ldots, y_{n}\right) \in B_{\ell_{q}^{n}}\right\}$, is equal to the following quantity up to lower order factors:

$$
\left\{\begin{array}{lll}
n^{\frac{3}{2}} & \text { if } & \max \{p, q\} \geqslant 2, \\
n^{1+\frac{1}{\max p, p q\}}} & \text { if } & \max \{p, q\} \leqslant 2
\end{array}\right.
$$

All of the results in Corollary 4 are new, except for the range $1 \leqslant p \leqslant 2$ of (3), which is due to [CCG ${ }^{+} 98$ ]. The range $p \in(2, \infty]$ of (3) is $\operatorname{SEP}\left(\ell_{p}^{n}\right)=\sqrt{n}$, which is incompatible with the statement $\operatorname{SEP}\left(\ell_{p}^{n}\right)=n^{1-1 / p}$ of [CCG ${ }^{+} 98$. We will explain the reason why the latter assertion of [CCG ${ }^{+} 98$ is erroneous in Remark 77 .

The wealth of knowledge that is available on the volumetric quantities that appear in (1) leads to new estimates that relate the separation modulus of an $n$-dimensional normed space $\mathbf{X}$ to classical invariants of $\mathbf{X}$. We will derive several such results herein, without attempting to be encyclopedic. As a noteworthy example, we will deduce from the first inequality in (1) that if $B_{\mathbf{X}}$ is a polytope with $\rho n$ vertices, then

$$
\begin{equation*}
\operatorname{SEP}(\mathbf{X}) \gtrsim \frac{n}{\sqrt{\log \rho}} . \tag{6}
\end{equation*}
$$

We will also deduce that if $T_{2}(\mathbf{X})$ denotes the type 2 constant of $\mathbf{X}$ (see (78) or the survey [Mau03]), then

$$
\begin{equation*}
\operatorname{SEP}(\mathbf{X}) \gtrsim \max \left\{\sqrt{\operatorname{dim}(\mathbf{X})}, T_{2}(\mathbf{X})^{2}\right\} . \tag{7}
\end{equation*}
$$

We will see that both (6) and (7) are sharp for the entire range of the relevant parameters (e.g. in the two extremes, the case $\mathbf{X}=\ell_{1}^{n}$ corresponds to $\rho=O(1)$ and $T_{2}(\mathbf{X})=\sqrt{n}$ in (6) and (7), respectively, and the case when $\mathbf{X}$ is $O(1)$-isomorphic to $\ell_{2}^{n}$ corresponds to $\log \rho=n$ and $T_{2}(\mathbf{X})=O(1)$ in (6) and (7), respectively).
1.1.1. A conjectural isomorphic reverse isoperimetric phenomenon. The lower bound on $\operatorname{SEP}(\mathbf{X})$ in Theorem 3 is not always sharp. Indeed, consider $\mathbf{X}=\ell_{1}^{n} \oplus \ell_{2}^{n}$ for which $\operatorname{SEP}(\mathbf{X})=n$ yet $\operatorname{vr}\left(\mathbf{X}^{*}\right) \sqrt{\operatorname{dim}(\mathbf{X})}=n^{3 / 4}$. It could be, however, that the upper bound on $\operatorname{SEP}(\mathbf{X})$ in Theorem 3 is optimal for every $\mathbf{X}$.

Question 5. Is it true that the separation modulus of any normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is bounded above and below by universal constant multiples of the minimum of $\operatorname{diam}_{\mathbf{X}^{*}}\left(\Pi B_{\mathbf{Y}}\right) / \mathrm{vol}_{n}\left(B_{\mathbf{Y}}\right)$ over all the normed spaces $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ that satisfy $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$ ?

See Remark 23 for an explanation why the minimum that is described in Question 5 is affine invariant, which is necessary for Question 5 to make sense, since the separation modulus is abi-Lipschitz invariant. For sufficiently symmetric spaces, we expect that the lower bound on $\operatorname{SEP}(\mathbf{X})$ in Theorem 3 is sharp.

Conjecture 6. Every finite dimensional normed space $\mathbf{X}$ with enough symmetries satisfies

$$
\begin{equation*}
\operatorname{SEP}(\mathbf{X})=\operatorname{vr}\left(\mathbf{X}^{*}\right) \sqrt{\operatorname{dim}(\mathbf{X})} . \tag{8}
\end{equation*}
$$

The notion of having enough symmetries was introduced in [GG71]; its definition is recalled in Section 1.6 .2 . We prefer to formulate Conjecture 6 using this notion at the present introductory juncture even though weaker requirements are needed for our purposes because it is a standard assumption in Banach space theory and it suffices for all of the most pressing applications that we have in mind.

The upper bound on $\operatorname{SEP}(\mathbf{X})$ in [8] implies by [LN05] that $\mathrm{e}(\mathbf{X}) \lesssim \operatorname{vr}\left(\mathbf{X}^{*}\right) \sqrt{\operatorname{dim}(\mathbf{X})}$, which would be a valuable Lipschitz extension theorem due to the fact that estimating the volume ratio is typically tractable given the variety of tools and extensive knowledge that are available in the literature. For example, Milman and Pisier [MP86] proved (improving by lower-order factors over a major theorem of Bourgain and Milman BM85, BM87]; see also (Mil87]), that any finite dimensional normed space $\mathbf{X}$ satisfies

$$
\begin{equation*}
\operatorname{vr}(\mathbf{X}) \lesssim C_{2}(\mathbf{X})\left(1+\log C_{2}(\mathbf{X})\right) \tag{9}
\end{equation*}
$$

where $C_{2}(\mathbf{X})$ is the cotype 2 constant of $\mathbf{X}$ (see (78) or the survey [Mau03]). Therefore, if (8] holds, then

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim C_{2}(\mathbf{X})\left(1+\log C_{2}(\mathbf{X})\right) \sqrt{\operatorname{dim}(\mathbf{X})} \tag{10}
\end{equation*}
$$

which would be a remarkable generalization of the bound $\mathrm{e}\left(\ell_{2}^{n}\right) \lesssim \sqrt{n}$ of LN05.
We expect that Theorem3already implies Conjecture 6, as expressed in the following conjecture which would yield a positive answer to Question 5 for normed spaces with enough symmetries.

Conjecture 7. If $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a normed space with enough symmetries, then there is a normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ that satisfies $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$ and $\operatorname{diam}_{\mathbf{X}^{*}}\left(\Pi B_{\mathbf{Y}}\right) / \operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right) \lesssim \operatorname{vr}\left(\mathbf{X}^{*}\right) \sqrt{n}$.

As an illustrative example of Conjecture 7 , consider $\mathbf{X}=\ell_{\infty}^{n}$. Then $\operatorname{vr}\left(\left(\ell_{\infty}^{n}\right)^{*}\right)=\operatorname{vr}\left(\ell_{1}^{n}\right)=O(1)$. One can compute that $\Pi B_{\ell_{\infty}^{n}}=2^{n-1} B_{\ell_{\infty}^{n}}$. Hence, $\operatorname{diam}_{\ell_{1}^{n}}\left(\Pi B_{\ell_{\infty}^{n}}\right) / \operatorname{vol}_{n}\left(B_{\ell_{\infty}^{n}}\right)=n$, so taking $\mathbf{Y}=\ell_{\infty}^{n}$ in Theorem 3 only gives the bound $\operatorname{SEP}\left(\ell_{\infty}^{n}\right) \lesssim n$. However, we will later see that there exists a normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ with $B_{\mathbf{Y}} \subseteq B_{\ell_{\infty}^{n}}$ for which $\operatorname{diam}_{\ell_{1}^{n}}\left(\Pi B_{\mathbf{Y}}\right) / \operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right) \lesssim \sqrt{n}$. More generally, we will prove that Conjecture 7 (hence also Conjecture 6, by Theorem 3 ) holds for any normed space for which the standard basis of $\mathbb{R}^{n}$ is 1 -symmetric, and we will also see that Conjecture 7 holds up to a logarithmic factor for its unitary ideal.

The minimization in Question 5 can be viewed as a shape optimization problem HP18 that could potentially be approached using calculus of variations. Given an origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$, it asks for the minimum of the affine invariant functional $L \mapsto$ outradius $_{K^{\circ}}(\Pi L) /$ vol $_{n}(L)$ over all originsymmetric convex bodies $L \subseteq K$, where for any two origin-symmetric convex bodies $A, B \subseteq \mathbb{R}^{n}$ we denote the minimum radius of a dilate of $A$ that circumscribes $B$ by outradius ${ }_{A}(B)=\min \{r \geqslant 0: B \subseteq r A\}$, and $K^{\circ}=\left\{y \in \mathbb{R}^{n}: \sup _{x \in K}\langle x, y\rangle \leqslant 1\right\}$ is the polar of $K$. Conjecture 7 asserts that if $K$ has enough symmetries, then this minimum is bounded above and below by universal constant multiples of $\operatorname{vr}\left(K^{\circ}\right) \sqrt{n}$.

The minimization problem in Question5also has an isoperimetric flavor. As such, its investigation led us to formulate the following conjectural twist of Ball's reverse isoperimetric phenomenon [Bal91c], which we think is a fundamental geometric open question and it would be valuable to understand it even without its consequences that we derive herein.

The isoperimetric quotient of a convex body $K \subseteq \mathbb{R}^{n}$ is defined (see Had57, page 269] or [Sch89]) to be

$$
\begin{equation*}
\operatorname{iq}(K)=\frac{\operatorname{vol}_{n-1}(\partial K)}{\operatorname{vol}_{n}(K)^{\frac{n-1}{n}}} . \tag{11}
\end{equation*}
$$

Using this notation, the classical Euclidean isoperimetric theorem states that

$$
\begin{equation*}
\operatorname{iq}(K) \geqslant \operatorname{iq}\left(B_{\ell_{2}^{n}}\right)=\frac{n \sqrt{\pi}}{\Gamma\left(\frac{n}{2}+1\right)^{\frac{1}{n}}}=\sqrt{n} \tag{12}
\end{equation*}
$$

The following theorem of Ball $[\overline{B a l 91 c}]$ shows that a judicious choice of the scalar product on $\mathbb{R}^{n}$ ensures that the isoperimetric quotient of a convex body can also be bounded from above.

Theorem 8 (Ball's reverse isoperimetric theorem Bal91c). For every $n \in \mathbb{N}$ and every origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$ there exists a linear transformation $S \in \mathrm{SL}_{n}(\mathbb{R})$ such that $\mathrm{iq}(S K) \leqslant 2 n=\mathrm{iq}\left([-1,1]^{n}\right)$.

We expect that in the isomorphic regime (i.e., permitting non-isometric $O(1)$ perturbations), originsymmetric convex bodies have asymptotically better reverse isoperimetric properties than what is guaranteed by Theorem 8 . In fact, we conjecture that if in addition to passing from $K$ to $S K$ for some $S \in S L_{n}(\mathbb{R})$, a $O(1)$-perturbation of $S K$ is allowed, then the isoperimetric quotient can be decreased to be of the same order of magnitude as that of the Euclidean ball.

Conjecture 9 (isomorphic reverse isoperimetry). There exists a universal constant $c>0$ with the following property. For every $n \in \mathbb{N}$ and every origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$, there exist a linear transformation $S \in S L_{n}(\mathbb{R})$ and an origin-symmetric convex body $L \subseteq \mathbb{R}^{n}$ with $c S K \subseteq L \subseteq S K$ and $\mathrm{iq}(L) \lesssim \sqrt{n}$.

Conjecture 9 can be restated analytically as the assertion that any $n$-dimensional normed space is at Banach-Mazur distance $O(1)$ from a normed space whose unit ball has isoperimetric quotient $O(\sqrt{n})$. We will prove that Conjecture 9 holds when $K$ is the unit ball of $\ell_{p}^{n}$ for any $p \in[1, \infty]$ and $n \in \mathbb{N}$, and we will also see that Conjecture 9 holds up to lower-order factors for any Schatten-von Neumann trace class.

The requirement $L \supseteq c S K$ of Conjecture 9 implies that $\sqrt[n]{\operatorname{vol}_{n}(L)} \geqslant c \sqrt[n]{\operatorname{vol}_{n}(K)}$. So, the following weaker conjecture is implied by Conjecture 9, we will prove it for any 1-unconditional body.

Conjecture 10 (weak isomorphic reverse isoperimetry). For every $n \in \mathbb{N}$ and every origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$ there exist a linear transformation $S \in S L_{n}(\mathbb{R})$ and an origin-symmetric convex body $L \subseteq S K$ that satisfies $\sqrt[n]{\operatorname{vol}_{n}(L)} \gtrsim \sqrt[n]{\operatorname{vol}_{n}(K)}$ and $\operatorname{iq}(L) \lesssim \sqrt{n}$.

In Section 1.6 we will elucidate the relation between the task of bounding from above the rightmost quantity in (3) and isomorphic reverse isoperimetry. While Conjecture 9 is the strongest version of the isomorphic reverse isoperimetric phenomenon that we expect holds in full generality, we will see that it would suffice to prove its weaker variant Conjecture 10 for the purpose of using Theorem 3 . In particular, consider the following symmetric version of Conjecture 10 which we will prove in Section 1.6 implies Conjecture 7 (hence, using Theorem 3 it also implies Conjecture 6).
Conjecture 11 (symmetric version of Conjecture 10 . For every $n \in \mathbb{N}$, if $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a normed space with enough symmetries whose isometry group is a subgroup of the orthogonal group $\mathrm{O}_{n} \subseteq \mathrm{GL}_{n}(\mathbb{R})$, then there is a normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ with $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$ and $\sqrt[n]{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \gtrsim \sqrt[n]{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}$ such that $\mathrm{iq}\left(B_{\mathbf{Y}}\right) \lesssim \sqrt{n}$.

The only difference between Conjecture 10 and Conjecture 11 is that if we impose the further requirement that $K$ is the unit ball of a normed space with enough symmetries whose isometry group consists only of orthogonal matrices, then we are naturally conjecturing that $S$ can be taken to be the identity matrix, i.e., there is no need to change the standard Euclidean structure on $\mathbb{R}^{n}$.

We will prove Conjecture 11 for various spaces, including $\ell_{p}^{n}\left(\ell_{q}^{n}\right)$ for any $p, q \geqslant 1$ and $n \in \mathbb{N}$, and any finite dimensional space with a 1 -symmetric basis. Also, we will show that Conjecture 11 holds up to a factor of $O(\sqrt{\log n})$ for any unitarily invariant norm on $\mathrm{M}_{n}(\mathbb{R})$. In general, an argument that was shown to us by B. Klartag and E. Milman and is included in Section 7 (see also Section 1.6.3) shows that Conjecture 10 and Conjecture 11 hold up to a factor of $O(\log n)$. We will see that these results lead to Corollary 4 , and in general we will deduce that Conjecture 7 , and hence, thanks to Theorem 3 , also Conjecture 6 , hold up to lower order factors. Thus, we will obtain the following theorem.

Theorem 12. $\operatorname{SEP}(\mathbf{X})=\operatorname{vr}\left(\mathbf{X}^{*}\right) \operatorname{dim}(\mathbf{X})^{\frac{1}{2}+o(1)}$ for any normed space $\mathbf{X}$ with enough symmetries.
Assuming Conjecture 11 it is possible to compute the exact asymptotic growth rate of the separation moduli of several important matrix spaces. For example, if Conjecture 11 holds for $S_{\infty}^{n}$, then we will see that the $o(1)$ term in (4) could be removed altogether, i.e.,

$$
\begin{equation*}
\forall(p, n) \in[1, \infty] \times \mathbb{N}, \quad \operatorname{SEP}\left(\mathrm{S}_{p}^{n}\right)=n^{\max \left\{1, \frac{1}{2}+\frac{1}{p}\right\}} . \tag{13}
\end{equation*}
$$

Also, assuming Conjecture 11 the lower order factors in the last two statements of Corollary 4 could be removed, namely we will see that Conjecture 11 implies that the separation modulus of $\mathrm{M}_{n}(\mathbb{R})$ equipped with the operator norm $\|\cdot\|_{\ell_{p}^{n} \rightarrow \ell_{q}^{n}}$ from $\ell_{p}^{n}$ to $\ell_{q}^{n}$ satisfies

$$
\operatorname{SEP}\left(\mathrm{M}_{n}(\mathbb{R}),\|\cdot\|_{\ell_{p}^{n} \rightarrow \ell_{q}^{n}}\right)= \begin{cases}n^{\frac{3}{2}-\frac{1}{\min \mid p, q]}} & \text { if } p, q \geqslant 2,  \tag{14}\\ n^{\frac{1}{2}+\frac{\max \mid p, q]}{m}} & \text { if } p, q \leqslant 2, \\ n & \text { if } p \leqslant 2 \leqslant q, \\ n^{\max \left\{1, \frac{1}{q}-\frac{1}{p}+\frac{1}{2}\right\}} & \text { if } q \leqslant 2 \leqslant p,\end{cases}
$$

and the separation modulus of the projective tensor product $\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}$ satisfies

$$
\operatorname{SEP}\left(\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}\right)=\left\{\begin{array}{lll}
n^{\frac{3}{2}} & \text { if } & \max \{p, q\} \geqslant 2  \tag{15}\\
n^{1+\frac{1}{\max }[p, q]} & \text { if } & \max \{p, q\} \leqslant 2
\end{array}\right.
$$

Remark 173 describes ramifications of these conjectural statements to norms of algorithmic importance.
Roadmap. The rest of the Introduction effectively restarts the description of the present work, with many more details/definitions/background/ideas of proofs, than what we have included above. We organized the introductory material in this way since this work pertains to multiple mathematical disciplines, including Banach spaces, convex geometry, nonlinear functional analysis, metric embeddings, extension of functions, and theoretical computer science. The backgrounds of potential readers are therefore varied, so even though the above overview achieves the goal of presenting the main results quickly, it inevitably includes terminology that is not familiar to some. The aforementioned organizational choice makes the ensuing discussion accessible. Additional background can be found in the monographs [LT77, MS86, TJ89] (Banach space theory), [BL00] (nonlinear functional analysis), [Mat02, Ost13] (metric embeddings), [BB12] (extension of functions), as well as the references that are cited throughout.

While the ensuing extended introductory text is not short, it achieves more than merely a description of the results, history, concepts and methods: It also contains groundwork that is needed for the subsequent sections. Thus, reading the Introduction will lead to a thorough conceptual understanding of the contents, leaving to the remaining sections considerations that are for the most part more technical.

We will start by focusing on the classical Lipschitz extension problem because it is more well known than the stochastic clustering issues that lead to most of our new results on Lipschitz extension, and also because it requires less technicalities (e.g. a suitable measurability setup) than our subsequent treatment of stochastic clustering. Throughout the Introduction (and beyond), we will formulate conjectures and questions that are valuable even without the links to Lipschitz extension and clustering that are derived herein. After the Introduction, the rest of this work will be organized thematically as follows. Section 2 is devoted to proofs of our various lower bounds, namely impossibility results that rule out the existence of extensions and clusterings with certain properties. Section 3 and Section 4 deal with positive results about random partitions. Specifically, Section 3 is of a more foundational nature as it describes the concepts, basic constructions, and proofs of measurability statements that are needed for later applications in the infinitary setting (of course, measurability can be ignored for statements about finite sets). Section 4 analyses in the case of normed spaces a periodic version of a commonly used randomized partitioning technique called iterative ball partitioning, and computes optimally (up to universal constant
factors) the probabilities of its separation and padding events. Section 5 shows how to pass from random partitions to Lipschitz extension, by adjusting to the present setting the method that was developed in [LN05]. Section 5 also contains further foundational results on Lipschitz extension, as well questions and conjectures that are of independent interest. Section 6 contains a range of volume and surface area estimates that are needed in conjunction with the theorems of the preceding sections in order to deduce new Lipschitz extension and stochastic clustering results for various normed spaces and their subsets. Section 7 proves that Conjecture 10 and Conjecture 11 hold up to a factor of $O(\log n)$, and also shows that the approach that leads to this result cannot fully resolve Conjecture 11.
1.2. Basic notation. Given a metric space ( $m, d_{m}$ ), a point $x \in T$ and a radius $r \geqslant 0$, the corresponding closed ball is denoted $B_{m}(x, r)=\left\{y \in m: d_{m}(y, x) \leqslant r\right\}$. If $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ is a Banach space (in this work, all vector spaces are over the real scalars unless stated otherwise), then denote by $B_{\mathbf{X}}$ the unit ball centered at the origin. Under this notation we have $B_{\mathbf{X}}=B_{\mathbf{X}}(0,1)$ and $B_{\mathbf{X}}(x, r)=x+r B_{\mathbf{X}}$ for every $x \in X$ and $r \geqslant 0$.

If $\left(m, d_{m}\right),\left(n, d_{n}\right)$ are metric spaces and $\psi: M \rightarrow n$, then for $\mathcal{C} \subseteq m$ the Lipschitz constant of $\psi$ on $\mathcal{C}$ is denoted $\|\psi\|_{\text {Lip }(\mathrm{e} ; n)} \in[0, \infty]$. Thus, if $\mathcal{C}$ contains at least two points, then

$$
\|\psi\|_{\text {Lip }(e ; \eta)} \stackrel{\text { def }}{=} \sup _{\substack{x, y \in \mathcal{C} \\ x \neq y}} \frac{d_{n}(\psi(x), \psi(y))}{d_{m}(x, y)}
$$

In the special case $n=\mathbb{R}$ we will use the simpler notation $\|\psi\|_{\text {Lip }(\mathcal{C} ; \mathbb{R})}=\|\psi\|_{\text {Lip }(\mathcal{C})}$.
If $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right),\left(\mathbf{Y},\|\cdot\|_{\mathbf{Y}}\right)$ are isomorphic Banach spaces, then their Banach-Mazur distance $d_{\mathrm{BM}}(\mathbf{X}, \mathbf{Y})$ is the infimum of the products of the operator norms $\|T\|_{\mathbf{X} \rightarrow \mathbf{Y}}$ and $\left\|T^{-1}\right\|_{\mathbf{Y} \rightarrow \mathbf{X}}$ over all possible linear isomorphisms $T: \mathbf{X} \rightarrow \mathbf{Y}$. The (bi-Lipschitz) distortion of a metric space $\left(m, d_{m}\right)$ into a metric space $\left(n, d_{n}\right)$, denoted $\mathrm{c}_{\left(n, d_{n}\right)}\left(m, d_{m}\right)$ or $\mathrm{c}_{n}(m)$ if the underlying metrics are clear from the context, is the infimum over those $D \in[1, \infty]$ for which there exists a mapping $\phi: m \rightarrow \eta$ and (a scaling factor) $\lambda>0$ such that

$$
\begin{equation*}
\forall x, y \in M, \quad \lambda d_{m}(x, y) \leqslant d_{n}(\phi(x), \phi(y)) \leqslant D \lambda d_{m}(x, y) \tag{16}
\end{equation*}
$$

Fix $n \in \mathbb{N}$. Throughout what follows, $\mathbb{R}^{n}$ will be always be endowed with its standard Euclidean structure, i.e., with the scalar product $\langle x, y\rangle=x_{1} y_{1}+\ldots+x_{n} y_{n}$ for $x=\left(x_{1}, \ldots, x_{n}\right), y=\left(y_{1}, \ldots, y_{n}\right) \in \mathbb{R}^{n}$. Given $z \in \mathbb{R}^{n} \backslash\{0\}$, the orthogonal projection onto its orthogonal hyperplane $z^{\perp}=\left\{x \in \mathbb{R}^{n}:\langle x, z\rangle=0\right\}$ will be denoted $\operatorname{Proj}_{z^{\perp}}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$. For $0<s \leqslant n$, the $s$-dimensional Hausdorff measure of a closed subset $A \subseteq \mathbb{R}^{n}$ is denoted $\operatorname{vol}_{s}(A)$. Integration with respect to the $s$-dimensional Hausdorff measure is indicated by $\mathrm{d} x$. If $0<\operatorname{vol}_{s}(A)<\infty$ and $f: A \rightarrow \mathbb{R}$ is continuous, then write $f_{S} f(x) \mathrm{d} x=\operatorname{vol}_{s}(A)^{-1} \int_{A} f(x) \mathrm{d} x$.

Given a normed space $(\mathbf{X},\|\cdot\| \mathbf{x})$ and $p \in[1, \infty], \ell_{p}^{n}(\mathbf{X})$ is the vector space $\mathbf{X}^{n}$ equipped with the norm

$$
\forall x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbf{X}^{n}, \quad\|x\|_{\ell_{p}^{n}(\mathbf{X})}=\left(\left\|x_{1}\right\|_{\mathbf{X}}+\ldots+\left\|x_{n}\right\|_{\mathbf{X}}\right)^{\frac{1}{p}},
$$

where for $p=\infty$ this is understood to be $\|x\|_{\ell_{\infty}^{n}(\mathbf{X})}=\max _{j \in\{1, \ldots, n\}}\left\|x_{j}\right\|_{\mathbf{X}}$. It is common to use the simpler notation $\ell_{p}^{n}=\ell_{p}^{n}(\mathbb{R})$ and we write as usual $S^{n-1}=\partial B_{\ell_{2}^{n}}$. The Schatten-von Neumann trace class $S_{p}^{n}$ is the ( $n^{2}$-dimensional) space of all $n$ by $n$ real matrices $\mathrm{M}_{n}(\mathbb{R})$, equipped with the norm that is defined by

$$
\forall T \in \mathrm{M}_{n}(\mathbb{R}), \quad\|T\|_{S_{p}^{n}}=\left(\operatorname{Tr}\left(\left(T T^{*}\right)^{\frac{p}{2}}\right)\right)^{\frac{1}{p}}=\left(\operatorname{Tr}\left(\left(T^{*} T\right)^{\frac{p}{2}}\right)\right)^{\frac{1}{p}}
$$

where $\|T\|_{S_{\infty}^{n}}=\|T\|_{\ell_{2}^{n} \rightarrow \ell_{2}^{n}}$ is the operator norm of $T$ when it is viewed as a linear operator from $\ell_{2}^{n}$ to $\ell_{2}^{n}$.
1.3. Lipschitz extension. As we recalled in Section 1.1 one associates to every metric space ( $m, d_{m}$ ) a bi-Lipschitz invarian $\sqrt{6}$ called the Lipschitz extension modulus of $\left(m, d_{m}\right)$ and denoted e( $m, d_{m}$ ) or $\mathrm{e}(M)$ if the metric is clear from the context, by defining it to be the infimum over those $K \in[1, \infty]$ with the property that for every nonempty subset $\mathcal{C} \subseteq m$, every Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ and every Lipschitz function $f: \mathcal{C} \rightarrow \mathbf{Z}$ there is a mapping $F: m \rightarrow \mathbf{Z}$ that extends $f$, i.e., $F(x)=f(x)$ whenever $x \in \mathcal{C}$, and

[^2]$\|F\|_{\text {Lip }(m, \mathbf{Z})} \leqslant K\|f\|_{\text {Lip }(\mathrm{e}, \mathbf{Z})}$; see Figure 1 . All of the ensuing extension theorems hold for a larger class of target metric spaces that need not necessarily be Banach spaces, including Hadamard spaces and Busemann nonpositively curved spaces [BH99], or more generally spaces that posses a conical geodesic bicombing (see e.g. [DL15]). This greater generality will be discussed in Section5, but we prefer at this introductory juncture to focus on the more classical and highly-studied setting of Banach space targets.


FIGURE 1. Given $K \geqslant 1$, the assertion that the Lipschitz extension modulus of a metric space $M$ satisfies $\mathrm{e}(M)<K$ means that for all subsets $\mathcal{C} \subseteq M$, all Banach spaces $\mathbf{Z}$ and all 1-Lipschitz mappings $f: \mathcal{C} \rightarrow \mathbf{Z}$, there is a $K$-Lipschitz mapping $F: M \rightarrow \mathbf{Z}$ such that the above diagram commutes, where $\mathrm{Id}_{\mathrm{C} \rightarrow m}: \mathcal{C} \rightarrow M$ is the formal inclusion.

When $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ is a finite dimensional normed space, the currently best-available general bounds on the quantity $\mathrm{e}(\mathbf{X})$ in terms of $\operatorname{dim}(\mathbf{X})$ are contained the following theorem.

Theorem 13. There is a universal constant $c>0$ such that for any finite dimensional normed space $\mathbf{X}$,

$$
\begin{equation*}
\operatorname{dim}(\mathbf{X})^{c} \lesssim \mathrm{e}(\mathbf{X}) \lesssim \operatorname{dim}(\mathbf{X}) . \tag{17}
\end{equation*}
$$

The bound $\mathrm{e}(\mathbf{X}) \lesssim \operatorname{dim}(\mathbf{X})$ in (17) is a famous result of Johnson, Lindenstrauss and Schechtman [JLS86], which they proved by cleverly refining the classical extension method of Whitney Whi34; different proofs of this estimate were found by Lee and the author [LN05] as well as by Brudnyi and Brudnyi [BB06] (see also the discussion in the paragraph following equation (37) below). It remains a major longstanding open problem to determine whether or not the bound of [JLS86] could be improved to $\mathrm{e}(\mathbf{X})=o(\operatorname{dim}(\mathbf{X})$ ).

The new content of Theorem 13 is the lower bound on $\mathrm{e}(\mathbf{X})$, which improves over the previously known bound $\mathrm{e}(\mathbf{X}) \geqslant \exp (c \sqrt{\log \operatorname{dim}(\mathbf{X})})$; see Remark 97 for the history of this question. It is a very interesting open problem to determine the supremum over those $c$ for which Theorem 13 holds. ${ }^{7}$ More generally, it is natural to aim to evaluate the precise power-type behavior of $\mathrm{e}(\mathbf{X})$ as $\operatorname{dim}(\mathbf{X}) \rightarrow \infty$ for specific (sequences of) finite dimensional normed spaces $\mathbf{X}$. However, prior to the present work and despite many efforts over the years, this was not achieved for any finite dimensional normed space whatsoever.

Theorem 14 (restatement of Theorem 2). For every $n \in \mathbb{N}$ we have $\mathrm{e}\left(\ell_{\infty}^{n}\right)=\sqrt{n}$.
The bound e $\left(\ell_{\infty}^{n}\right) \gtrsim \sqrt{n}$ follows from a combination of [BB05, Theorem 4] and [BB07a, Theorem 1.2]. The new content of Theorem 14 is the the upper bound $\mathrm{e}\left(\ell_{\infty}^{n}\right) \lesssim \sqrt{n}$ (and, importantly, the extension procedure that leads to it; see below). The previously best-known upper bound on $\mathrm{e}\left(\ell_{\infty}^{n}\right)$ was the aforementioned $O(n)$ estimate of [JLS86]. The question of evaluating the asymptotic behavior of $\mathrm{e}\left(\ell_{p}^{n}\right)$ as $n \rightarrow \infty$ for each $p \in[1, \infty]$ is natural and longstanding; it was stated in [BB05] Problem 2] and reiterated in [BB07b, Section 4], BB07a, Problem 1.4] and [BB12, Problem 8.14]. Theorem 14 answers this question when $p=\infty$. The upper bound on $\mathrm{e}\left(\ell_{\infty}^{n}\right)$ of Theorem 14 is a special case of a general extension criterion that provides the best-known Lipschitz extension results in other settings (including for $\ell_{p}^{n}$ when $p>2$ ), but we chose to state it separately because it yields the first (and currently essentially only) family of normed spaces for which the growth rate of their Lipschitz extension moduli has been determined.

[^3]Remark 15. It is meaningful to study extension of $\theta$-Hölder functions for any $0<\theta \leqslant 1$. Namely, one can analogously define the $\theta$-Hölder extension modulus of a metric space $(m, d m)$, denoted $\mathrm{e}^{\theta}(m)$. Alternatively, this notion falls into the above Lipschitz-extension framework because one can define

$$
\begin{equation*}
\mathrm{e}^{\theta}(m) \stackrel{\text { def }}{=} \mathrm{e}\left(m, d_{m}^{\theta}\right) \tag{18}
\end{equation*}
$$

The results that we obtain herein also yield improved estimates on $\theta$-Hölder extension moduli; see Corollary 139 . However, when $\theta<1$ we never get a matching lower bound (the reason why we can do better in the Lipschitz regime $\theta=1$ is essentially due to the fact that Lipschitz functions are differentiable almost everywhere). For example, in the setting of Theorem 14 we get the upper bound

$$
\begin{equation*}
\forall \theta \in(0,1], \quad \mathrm{e}^{\theta}\left(\ell_{\infty}^{n}\right) \lesssim n^{\frac{\theta}{2}} \tag{19}
\end{equation*}
$$

but the best lower bound on $\mathrm{e}^{\theta}\left(\ell_{\infty}^{n}\right)$ that we are at present able to prove is

$$
\mathrm{e}^{\theta}\left(\ell_{\infty}^{n}\right) \gtrsim n^{\max \left\{\frac{\theta}{4}, \frac{\theta}{2}+\theta^{2}-1\right\}}=\left\{\begin{array}{lll}
n^{\frac{\theta}{4}} & \text { if } & 0 \leqslant \theta \leqslant \frac{\sqrt{65}-1}{8}  \tag{20}\\
n^{\frac{\theta}{2}+\theta^{2}-1} & \text { if } \quad \frac{\sqrt{65}-1}{8} \leqslant \theta \leqslant 1
\end{array}\right.
$$

We conjecture that $\mathrm{e}^{\theta}\left(\ell_{\infty}^{n}\right) \asymp_{\theta} n^{\frac{\theta}{2}}$, but proving this for $\theta<1$ would likely require a genuinely new idea.
Question 16. Despite its utility in many cases, the extension method that underlies Theorem 14 does not yield improved bounds for some important spaces, including notably $\ell_{1}^{n}$ and $\ell_{2}^{n}$. Thus, determining the asymptotic behavior of $\mathrm{e}\left(\ell_{1}^{n}\right)$ and $\mathrm{e}\left(\ell_{2}^{n}\right)$ as $n \rightarrow \infty$ remains a tantalizing open question. Specifically, the currently best-known bounds on $\mathrm{e}\left(\ell_{1}^{n}\right)$ are

$$
\begin{equation*}
\sqrt{n} \lesssim \mathrm{e}\left(\ell_{1}^{n}\right) \lesssim n \tag{21}
\end{equation*}
$$

where the first inequality in (21] is due to Johnson and Lindenstrauss [JL84] and the second inequality in (21) is the aforementioned general upper bound of [JLS86] on the Lipschitz extension modulus of any $n$-dimensional normed space. The currently best-known bounds in the Hilbertian setting are

$$
\begin{equation*}
\sqrt[4]{n} \lesssim \mathrm{e}\left(\ell_{2}^{n}\right) \lesssim \sqrt{n} \tag{22}
\end{equation*}
$$

where the first inequality in (22) is due to Mendel and the author [MN13] (a different proof of this lower bound on $\mathrm{e}\left(\ell_{2}^{n}\right)$ follows from [Nao21b]), and the second inequality in (22) is from [LN05].

By the bi-Lipschitz invariance of the Lipschitz extension modulus, the second inequality in (22) implies the following bound from [LN05], which holds for every finite dimensional normed space $\mathbf{X}$.

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim d_{\mathrm{BM}}\left(\mathbf{X}, \ell_{2}^{\operatorname{dim}(\mathbf{X})}\right) \sqrt{\operatorname{dim}(\mathbf{X})} \tag{23}
\end{equation*}
$$

This refines the upper bound on $\mathrm{e}(\mathbf{X})$ in (17) because $d_{\mathrm{BM}}\left(\mathbf{X}, \ell_{2}^{\operatorname{dim}(\mathbf{X})}\right) \leqslant \sqrt{\operatorname{dim}(\mathbf{X})}$ by John's theorem (Joh48].
Remark 17. In the context of the aforementioned question if the bound $e(\mathbf{X}) \lesssim \operatorname{dim}(\mathbf{X})$ of [JLS86] is optimal, by (23) we see that $e(\mathbf{X})=o(\operatorname{dim}(\mathbf{X}))$ unless the Banach-Mazur distance between $\mathbf{X}$ and Euclidean space is of order $\sqrt{\operatorname{dim(X)}}$. Structural properties of such spaces of extremal distance to Euclidean space have been studied in [MW78, Pis79, Bou82, JS82a, ATTJ05]; see also chapters 6 and 7 of [TJ89]. In particular, the Mil'man-Wolfson theorem [MW78] asserts that this holds if and only if $\mathbf{X}$ has a subspace of dimension $k=k(\operatorname{dim}(\mathbf{X}))$ whose Banach-Mazur distance to $\ell_{1}^{k}$ is $O(1)$, where $\lim _{n \rightarrow \infty} k(n)=\infty$.

As $d_{\mathrm{BM}}\left(\ell_{p}^{n}, \ell_{2}^{n}\right)=n^{|p-2| /(2 p)}$ for all $n \in \mathbb{N}$ and $p \in[1, \infty]$ (see (JL01) Section 8]), it follows from (23) that

$$
\mathrm{e}\left(\ell_{p}^{n}\right) \lesssim \begin{cases}n^{\frac{1}{p}} & \text { if } p \in[1,2]  \tag{24}\\ n^{1-\frac{1}{p}} & \text { if } p \in[2, \infty]\end{cases}
$$

(24) was the previously best-known upper bound on $\mathrm{e}\left(\ell_{p}^{n}\right)$, and here we improve it for every $p>2$.

Theorem 18. For every $n \in \mathbb{N}$ and every $p \in[1, \infty]$ we have $\mathrm{e}\left(\ell_{p}^{n}\right) \lesssim n^{\max \left\{\frac{1}{2}, \frac{1}{p}\right\}}$.

Theorem 14 is the case $p=\infty$ of Theorem 18. We do not know if Theorem 18 is optimal (perhaps up to lower order factors) as $n \rightarrow \infty$ for fixed $p \in[2, \infty)$, but we conjecture that this is indeed the case, which would resolve [BB05, Problem 2]. The currently best-known lower bound on e( $\ell_{p}^{n}$ ) for every $p \in[1, \infty]$ is

$$
\mathrm{e}\left(\ell_{p}^{n}\right) \gtrsim \begin{cases}n^{\frac{1}{p}-\frac{1}{2}} & \text { if } 1 \leqslant p \leqslant \frac{4}{3}  \tag{25}\\ \sqrt[4]{n} & \text { if } \frac{4}{3} \leqslant p \leqslant 2 \\ n^{\frac{1}{2 p}} & \text { if } 2 \leqslant p \leqslant 3 \\ n^{\frac{1}{2}-\frac{1}{p}} & \text { if } 3 \leqslant p \leqslant \infty\end{cases}
$$

A lower bound on $\mathrm{e}\left(\ell_{p}^{n}\right)$ that coincides with (25) when $p \in[1,4 / 3] \cup[3, \infty]$ is stated in Corollary 8.12 of [BB12], but [BB12, Corollary 8.12] is weaker than (25) when $4 / 3<p<3$. The reason for this is that the lower bound of [MN13] on e $\left(\ell_{2}^{n}\right)$ that appears in (22) was not available when [BB12] was written, but (25) for $4 / 3<p<3$ follows quickly by combining the first inequality in (22) with [FLM77]; see Remark 2.4.

Remark 19. Theorem 18 resolves negatively a conjecture that A. Brudnyi and Y. Brudnyi posed as Conjecture 5 in [BB05]. They conducted a comprehensive study of the linear extension problem for real-valued Lipschitz functions, where one considers for a metric space $\left(m, d_{m}\right)$ a quantity $\lambda(m)$ which is defined the same as $\mathrm{e}(M)$, but with the further requirements that the function $f$ is real-valued and that the extended function $F$ depends linearly on $f$. Namely, $\lambda(m)$ is the infimum over those $K \in[1, \infty]$ such that for every $\mathcal{C} \subseteq m$ there is a linear operator $\operatorname{Ext}{ }_{\mathcal{C}}: \operatorname{Lip}(\mathcal{C}) \rightarrow \operatorname{Lip}(M)$ that assigns to every Lipschitz function $f: \mathcal{C} \rightarrow \mathbb{R}$ a function $\operatorname{Ext}_{e} f: m \rightarrow \mathbb{R}$ satisfying $\operatorname{Ext}_{e} f(s)=f(s)$ for every $s \in \mathcal{C}$, and

$$
\left\|\operatorname{Ext}_{\mathrm{e}} f\right\|_{\operatorname{Lip}(m)} \leqslant K\|f\|_{\operatorname{Lip}(\mathrm{e})} .
$$

They also considered a natural variant of this quantity when $m=\mathbf{X}$ is a Banach space, denoted $\lambda_{\text {conv }}(\mathbf{X})$, which is defined almost identically to $\lambda(\mathbf{X})$ except that now the subset $\mathcal{C}$ is only allowed to be any convex subset of $\mathbf{X}$ rather than a subset of $\mathbf{X}$ without any additional restriction. Conjecture 5 in [BB05] states that

$$
\begin{equation*}
\forall(p, n) \in[1, \infty] \times \mathbb{N}, \quad \lambda\left(\ell_{p}^{n}\right)={ }_{p} \lambda_{\mathrm{conv}}\left(\ell_{p}^{n}\right) \sqrt{n} . \tag{26}
\end{equation*}
$$

Theorem 18 implies that this conjecture is false for every $p \in(2, \infty]$. Indeed, the asymptotic behavior of $\lambda_{\text {conv }}\left(\ell_{p}^{n}\right)$ was evaluated in BB07b , Theorem 2.19], where it was shown that

$$
\forall p \in[1, \infty], \quad \lambda_{\operatorname{conv}}\left(\ell_{p}^{n}\right)=n^{\left|\frac{1}{2}-\frac{1}{p}\right|} .
$$

Consequently, $\lambda_{\text {conv }}\left(\ell_{p}^{n}\right) \sqrt{n}=n^{1-\frac{1}{p}}$ when $p>2$. Next, in BB07a] a quantity $v(m)$ was associated to a metric space $\left(m, d_{m}\right)$ by defining it almost identically to the definition of e( $M$ ), except that the target Banach space $\mathbf{Z}$ is allowed to be any finite dimensional Banach space rather than any Banach space whatsoever. By definition $v(m) \leqslant \mathrm{e}\left(m_{)}\right.$, but actually $\lambda(m)=v(m)$ thanks to [BB07a, Theorem 1.2] (see the work AP20 of Ambrosio and Puglisi for more on this "linearization phenomenon"). Using these results in combination with Theorem 18 , we see that for every $p \in(2, \infty]$, as $n \rightarrow \infty$ we have

$$
\lambda\left(\ell_{p}^{n}\right)=v\left(\ell_{p}^{n}\right) \leqslant \mathrm{e}\left(\ell_{p}^{n}\right) \lesssim \sqrt{n}=o\left(n^{1-\frac{1}{p}}\right) .
$$

Thus, $\lambda\left(\ell_{p}^{n}\right)=o\left(\lambda_{\operatorname{conv}}\left(\ell_{p}^{n}\right) \sqrt{n}\right)$ as $n \rightarrow \infty$ for any $p>2$, in contrast to the conjecture (26] of [BB05].
Prior to passing to the general Lipschitz extension theorem that underlies the new results that were described above, we will further illustrate its utility by stating one more concrete application. For each $p \in[1, \infty]$ and $n \in \mathbb{N}$, if $k \in\{1, \ldots, n\}$, then let $\left(\ell_{p}^{n}\right)_{\leqslant k}$ denote the subset of $\mathbb{R}^{n}$ consisting of those vectors with at most $k$ nonzero coordinates, equipped with the metric that is inherited from $\ell_{p}^{n}$.
Theorem 20. For every $p \in[1, \infty]$, every $n \in \mathbb{N}$ and every $k \in\{1, \ldots, n\}$ we have $\mathrm{e}\left(\left(\ell_{p}^{n}\right)_{\leqslant k}\right) \lesssim k^{\max \left\{\frac{1}{p}, \frac{1}{2}\right\}}$.

Theorem 18 is the special case $k=n$ and $p \geqslant 2$ of Theorem 20 . If $1 \leqslant p \leqslant 2$ and $k=n$, then Theorem 20 is the estimate (24), which is the best-known upper bound on $\mathrm{e}\left(\ell_{p}^{n}\right)$ for $p$ in this range. However, for general $k \in\{1, \ldots, n\}$ Theorem 20 yields a refinement of 24 in the entire range $p \in[1, \infty]$ which does not seem to follow from previously known results. In particular, the case $p=2$ of Theorem 20 becomes

$$
\begin{equation*}
\mathrm{e}\left(\left(\ell_{2}^{n}\right)_{\leqslant k}\right) \lesssim \sqrt{k} . \tag{27}
\end{equation*}
$$

Even though 27) concerns a Euclidean setting, its proof relies on a construction that employs a multiscale partitioning scheme using balls of an auxiliary metric on $\mathbb{R}^{n}$ that differs from the ambient Euclidean metric. The utility of such a non-Euclidean geometric reasoning despite the Euclidean nature of the question being studied is discussed further in Section 1.4.
1.4. A volumetric upper bound on the Lipschitz extension modulus. We will prove that Theorem 20 (hence also its special cases Theorem 14 and Theorem 18 is a consequence of Theorem 21 below, which is a Lipschitz extension theorem for subsets of finite dimensional normed spaces in terms of volumes of hyperplane projections of their unit balls. Throughout what follows, for dealing with volumetric notions we will adhere to the following conventions. Given $n \in \mathbb{N}$, when we say that $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a normed space we mean that the underlying vector space is $\mathbb{R}^{n}$ and that $\|\cdot\|_{\mathrm{X}}: \mathbb{R}^{n} \rightarrow[0, \infty)$ is a norm on $\mathbb{R}^{n}$. This is, of course, always achievable by fixing any scalar product on an $n$-dimensional normed space. While the ensuing statements hold in this setting, i.e., for an arbitrary identification of $\mathbf{X}$ with $\mathbb{R}^{n}$, a judicious choice of such an identification is beneficial; the discussion of this important matter is postponed to Section 1.6.2 because it is not needed for the initial description of the main results. We will continue using the notation $B_{\mathbf{X}}=\left\{x \in \mathbb{R}^{n}:\|x\|_{\mathbf{X}} \leqslant 1\right\}$ for the unit ball of $\mathbf{X}$. Also, given $\mathcal{C} \subseteq \mathbb{R}^{n}$ we denote by $\mathcal{C}_{\mathbf{X}}$ the metric space consisting of the set $\mathcal{C}$ equipped with the metric that is inherited from $\|\cdot\|_{\mathbf{X}}$. This notation is important for us because we will crucially need to simultaneously consider more than one norm on $\mathbb{R}^{n}$.

Theorem 21. Suppose that $n \in \mathbb{N}$ and that $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ and $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ are two normed spaces. Then, for every $\mathcal{C} \subseteq \mathbb{R}^{n}$ we have the following upper bound on the Lipschitz extension modulus of $\mathcal{C}_{\mathrm{X}}$.

$$
\begin{equation*}
\mathrm{e}\left(\mathrm{C}_{\mathbf{X}}\right) \lesssim\left(\sup _{\substack{x, y \in \mathcal{C} \\ x \neq y}} \frac{\|x-y\|_{\mathbf{X}}}{\|x-y\|_{\mathbf{Y}}}\right) \sup _{\substack{x, y \in \mathcal{C} \\ x \neq y}}\left(\frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{(x-y)^{\perp} B \mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \cdot \frac{\|x-y\|_{\ell_{2}^{n}}}{\|x-y\|_{\mathbf{X}}}\right) \tag{28}
\end{equation*}
$$

We will next discuss the geometric meaning of Theorem 21 and derive some of its consequences, including Theorem 20. Firstly, by homogeneity the case $\mathcal{C}=\mathbb{R}^{n}$ of (28) becomes

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim\left(\sup _{y \in \partial B_{\mathbf{Y}}}\|y\|_{\mathbf{X}}\right) \sup _{x \in \partial B_{\mathbf{X}}}\left(\frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}} B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)}\|x\|_{\ell_{2}^{n}}\right) \tag{29}
\end{equation*}
$$

The quantity $\sup _{y \in \partial B_{\mathbf{Y}}}\|y\|_{\mathbf{X}}$ in (29) is the norm $\left\|\mathrm{Id}_{n}\right\|_{\mathbf{Y} \rightarrow \mathbf{X}}$ of the identity matrix $\operatorname{ld}_{n} \in \mathrm{M}_{n}(\mathbb{R})$ as an operator from $\mathbf{Y}$ to $\mathbf{X}$. Alternatively, $\sup _{y \in \partial B_{\mathbf{Y}}}\|y\|_{\mathbf{X}}=\operatorname{diam}_{\mathbf{X}}\left(B_{\mathbf{Y}}\right) / 2$, where for each $\mathcal{C} \subseteq \mathbb{R}^{n}$ we denote its diameter with respect to the metric that $\mathbf{X}$ induces $\operatorname{bybam}_{\mathbf{X}}(\mathcal{C})=\sup _{x, y \in \mathcal{C}}\|x-y\|_{\mathbf{X}}$.

Given a convex body $K \subseteq \mathbb{R}^{n}$, let $\Pi^{*} K \subseteq \mathbb{R}^{n}$ be the polar of the projection body of $K$, which is defined to be the unit ball of the norm $\|\cdot\|_{\Pi^{*} K}$ on $\mathbb{R}^{n}$ that is given by setting

$$
\begin{equation*}
\forall x \in \mathbb{R}^{n} \backslash\{0\}, \quad\|x\|_{\Pi^{*} K} \stackrel{\text { def }}{=} \frac{1}{2} \int_{\partial K}\left|\left\langle x, N_{K}(y)\right\rangle\right| \mathrm{d} y=\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}} K\right)\|x\|_{\ell_{2}^{n}}, \tag{30}
\end{equation*}
$$

where $N_{K}(y) \in S^{n-1}$ denotes the unit outer normal to $\partial K$ at $y \in \partial K$ (which is uniquely defined almost everywhere with respect to the surface-area measure on $\partial K$ ), and the final equality in (30) is the Cauchy projection formula (see e.g. [Gar06, Appendix A]). The projection body $\Pi К$ of $K$ is the polar of $\Pi^{*} K$. These important notions were introduced by Petty Pet67]. When $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a normed space let $\Pi^{*} \mathbf{X}$ be the normed space whose unit ball is $\Pi^{*} B_{\mathbf{X}}$. Let $\Pi \mathbf{X}=\left(\Pi^{*} \mathbf{X}\right)^{*}$ be the normed space whose unit ball is $\Pi B_{\mathbf{X}}$.

By substituting (30) into (29) we get the following interpretation of our bound on e(X) in terms of analytic and geometric properties of projection bodies; it is worthwhile to state it as a separate corollary even
though it is only a matter of notation because of its intrinsic interest and also because these alternative viewpoints were useful for guiding some of the subsequent considerations.

Corollary 22. Any two normed spaces $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right), \mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ satisfy

$$
\begin{align*}
\mathrm{e}(\mathbf{X}) & \lesssim \frac{\operatorname{diam}_{\mathbf{X}}\left(B_{\mathbf{Y}}\right) \operatorname{diam}_{\Pi^{*} \mathbf{Y}}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)}=\frac{\left\|\mathrm{Id}_{n}\right\|_{\mathbf{Y} \rightarrow \mathbf{X}}\left\|\mathrm{Id}_{n}\right\|_{\mathbf{X} \rightarrow \Pi^{*} \mathbf{Y}}}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \\
& =\frac{\left\|\mathrm{Id}_{n}\right\|_{\mathbf{X} \rightarrow \mathbf{Y}}\left\|\mathrm{Id}_{n}\right\|_{\Pi \mathbf{Y} \rightarrow \mathbf{X}^{*}}}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)}=\frac{\operatorname{diam}_{\mathbf{X}}\left(B_{\mathbf{Y}}\right) \operatorname{diam}_{\mathbf{X}^{*}\left(\Pi B_{\mathbf{Y}}\right)}^{\operatorname{vol}\left(B_{\mathbf{Y}}\right)}}{} . \tag{31}
\end{align*}
$$

The penultimate step in (31) is duality (the norm of an operator equals the norm of its adjoint) and the final quantity in (31) relates Theorem 21 to the second estimate in Theorem3.

Remark 23. Corollary 22 has the right affine invariance. For $S \in S L_{n}(\mathbb{R})$ let $S \mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{S X}\right)$ be the normed space whose unit ball is $S B_{\mathbf{X}}$; equivalently, $\|x\|_{S \mathbf{X}}=\left\|S^{-1} x\right\|_{\mathbf{X}}$ for every $x \in \mathbb{R}^{n}$. Then $\mathbf{X}$ and $S \mathbf{X}$ are isometric as metric spaces, so $\mathrm{e}(S \mathbf{X})=\mathrm{e}(\mathbf{X})$. We have $(S \mathbf{X})^{*}=\left(S^{*}\right)^{-1} \mathbf{X}^{*}$ (by definition), and $\Pi\left(S B_{\mathbf{Y}}\right)=\left(S^{*}\right)^{-1} \Pi B_{\mathbf{Y}}$ by [Pet67]. From this we see that $\operatorname{diam}_{(S X)^{*}}\left(\Pi B_{S \mathbf{Y}}\right)=\operatorname{diam}_{\mathbf{X}^{*}}\left(\Pi B_{\mathbf{Y}}\right)$. Thus, the minimum of the right hand side of (31) over all normed spaces $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ is also invariant under the action of $S L_{n}(\mathbb{R})$.

The special case of Theorem 21 in which the normed space $Y$ coincides with the given normed space $\mathbf{X}$ is in itself a nontrivial bound on the Lipschitz extension modulus. Examining this special case first will help elucidate how the idea arose to introduce an auxiliary space $\mathbf{Y}$ that may differ from $\mathbf{X}$, and why this can yield stronger estimates. If $\mathbf{X}=\mathbf{Y}$, then the bound (28) becomes

$$
\begin{equation*}
\mathrm{e}\left(\bigodot_{\mathbf{X}}\right) \lesssim \sup _{\substack{x, y \in \mathcal{C} \\ x \neq y}}\left(\frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{(x-y)^{\perp}} B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} \cdot \frac{\|x-y\|_{\ell_{2}^{n}}}{\|x-y\|_{\mathbf{X}}}\right) \tag{32}
\end{equation*}
$$

Correspondingly, the bound 29 becomes

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim \sup _{z \in \partial B_{\mathbf{X}}}\left(\frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z^{\perp}} B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}\|z\|_{\ell_{2}^{n}}\right)=\frac{\operatorname{diam}_{\Pi^{*} \mathbf{X}}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} \tag{33}
\end{equation*}
$$

Even these weaker estimates suffice to obtain new results, e.g. we will see that this is so if $2 \leqslant p=O(1)$ and $\mathbf{X}=\ell_{p}^{n}$. However, as we will soon explain, (33) does not imply an upper bound on $\ell_{\infty}^{n}$ that is better than the aforementioned general bound of [JLS86]. Despite this shortcoming of (32) and (33) relative to (28), it is worthwhile to state these special cases of Theorem 21 separately because they are simpler than (28) and hence perhaps somewhat easier to remember. Moreover, a naïve way to enhance the applicability of (32) is to leverage the fact that the Lipschitz extension modulus is a bi-Lipschitz invariant, so that

$$
\mathrm{e}\left(\mathcal{C}_{\mathbf{X}}\right) \leqslant\left\|\mathrm{Id}_{n}\right\|_{\operatorname{Lip}\left(\mathcal{C}_{\mathbf{Y}}, \mathcal{C}_{\mathbf{X}}\right)}\left\|\mathrm{Id}_{n}\right\|_{\operatorname{Lip}\left(\mathcal{C}_{\mathbf{X}}, \mathcal{C}_{\mathbf{Y}}\right)} \mathrm{e}\left(\mathcal{C}_{\mathbf{Y}}\right)
$$

Consequently, by estimating e( $\left(_{Y}\right)$ through (32) we formally deduce from (32) that

$$
\begin{equation*}
\mathrm{e}\left(\mathcal{C}_{\mathbf{X}}\right) \lesssim\left(\sup _{\substack{x, y \in \mathcal{C} \\ x \neq y}} \frac{\|x-y\|_{\mathbf{X}}}{\|x-y\|_{\mathbf{Y}}}\right)\left(\sup _{\substack{x, y \in \mathcal{C} \\ x \neq y}} \frac{\|x-y\|_{\mathbf{Y}}}{\|x-y\|_{\mathbf{X}}}\right) \cdot \sup _{\substack{x, y \in \mathcal{C} \\ x \neq y}}\left(\frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{\left.(x-y)^{\perp} B_{\mathbf{Y}}\right)}^{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \cdot \frac{\|x-y\|_{\ell_{2}^{n}}}{\|x-y\|_{\mathbf{Y}}}\right) . . . . . . . .}{}\right. \tag{34}
\end{equation*}
$$

We do not see how to deduce Theorem 18 and Theorem 20 from (34). However, we will show that (34) suffices for proving Theorem 14 (as well as some other results that will be presented later). In summary, even the case of Theorem 21 in which the auxiliary space $\mathbf{Y}$ coincides with $\mathbf{X}$ is valuable, but Theorem 21 does not follow from merely combining its special case $\mathbf{Y}=\mathbf{X}$ with bi-Lipschitz invariance.

Given a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ and $z \in \mathbb{R}^{n} \backslash\{0\}$, the quantity

$$
\begin{equation*}
\frac{1}{n} \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z^{\perp}} B \mathbf{X}\right)\|z\|_{\ell_{2}^{n}} \tag{35}
\end{equation*}
$$

is equal to the volume of the cone

$$
\begin{equation*}
\operatorname{Cone}_{z}\left(B_{\mathbf{X}}\right) \stackrel{\text { def }}{=} \operatorname{conv}\left(\{z\} \cup \operatorname{Proj}_{z^{\perp}} B_{\mathbf{X}}\right) \subseteq \mathbb{R}^{n} \tag{36}
\end{equation*}
$$

whose base is the ( $n-1$ )-dimensional convex set $\operatorname{Proj}_{z^{\perp}} B_{\mathbf{X}} \subseteq z^{\perp}$ and whose apex is $z$. In (36) and throughout what follows, $\operatorname{conv}(\cdot)$ denotes the convex hull. Thus, the estimate (33) can be restated as follows.

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim n \sup _{z \in \partial B_{\mathbf{X}}} \frac{\operatorname{vol}_{n}\left(\operatorname{Cone}_{z}\left(B_{\mathbf{X}}\right)\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} . \tag{37}
\end{equation*}
$$

Through (37) we see that the geometric interpretation of the "bad spaces" $\mathbf{X}$ for (33) is that these are the spaces that have a "pointy direction" $z \in \partial B_{\mathbf{X}}$ for which the volume of the cone Cone $z_{z}\left(B_{\mathbf{X}}\right)$ is a significant fraction of the volume of $B_{\mathbf{X}}$. Examples will be presented next, but note first that a short geometric argument (see the proof of [GNS12] Lemma 5.1]) shows that $\operatorname{vol}_{n}\left(\operatorname{Cone}_{z}\left(B_{\mathbf{X}}\right)\right) \leqslant \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) / 2$, so the right hand side of (37) is at most $n / 2$. Hence, (33) is a refinement of the classical bound $\mathrm{e}(\mathbf{X}) \lesssim n$ of (JLS86].

Nevertheless, a "vanilla" application of (33) does not yield an asymptotically better estimate than that of (JLS86] even when $\mathbf{X}=\ell_{\infty}^{n}$. Indeed, $B_{\ell_{\infty}^{n}}=[-1,1]^{n}$ and a simple argument (see [CF86]) shows that

$$
\begin{equation*}
\forall z \in \mathbb{R}^{n} \backslash\{0\}, \quad \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z^{\perp}}[-1,1]^{n}\right)}{\operatorname{vol}_{n}\left([-1,1]^{n}\right)}=\frac{\|z\|_{\ell_{1}^{n}}}{2\|z\|_{2}^{n}} . \tag{38}
\end{equation*}
$$

So, by considering the all l's vector $z=\mathbf{1}_{\{1, \ldots, n\}} \in \partial B_{\ell_{\infty}^{n}}$ we see that for $\mathbf{X}=\ell_{\infty}^{n}$ the right hand side of (33) is at least $n / 2$. The right hand side of (33) is at least $n / 2$ when $\mathbf{X}=\ell_{1}^{n}$, as seen by taking $z=(1,0, \ldots, 0) \in \partial B_{\ell_{1}^{n}}$. Such "problematic" directions $z \in \partial B_{\mathbf{X}}$ can sometimes be the overwhelming majority of $\partial B_{\mathbf{X}}$. Consider Ball's counterexample [Bal91b] to the Shepard Problem [She64], which states that for any $n \in \mathbb{N}$ there is a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ such that $\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)=1$ yet $\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z^{\perp}} B_{\mathbf{X}}\right) \gtrsim \sqrt{n}$ for every $z \in S^{n-1}$. Since $\operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right) \leqslant(3 / \sqrt{n})^{n}$ while $\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)=1$, the proportion of those $z \in \partial B_{\mathbf{X}}$ for which $\|z\|_{\ell_{2}^{n}} \geqslant \sqrt{n} / 4$ tends to 1 as $n \rightarrow \infty$ (exponentially fast). Any such $z$ satisfies $\|z\|_{\ell_{2}^{n} \operatorname{vol}_{n-1}}\left(\operatorname{Proj}_{z^{\perp}} B_{\mathbf{X}}\right) / \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) \gtrsim n$.

These obstacles can sometimes be overcome by perturbing the given normed space $\mathbf{X}$ prior to invoking (33), i.e., by using of Theorem 21 with a suitably chosen auxiliary normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$. In particular, since by Hölder's inequality $\|\cdot\|_{\ell_{2}^{n}} \leqslant n^{1 / 2-1 / p}\|\cdot\|_{\ell_{p}^{n}}$ when $p \geqslant 2$, Theorem 18 follows from a substitution of the space $\mathbf{Y}_{p}^{n}$ of Theorem 24 below into Theorem 21 (with $\mathbf{X}=\ell_{p}^{n}$ ), or even into (34).

Theorem 24. For any $n \in \mathbb{N}$ and $p \in[1, \infty]$ there is a normed space $\mathbf{Y}_{p}^{n}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}_{p}^{n}}\right)$ that satisfies

$$
\begin{equation*}
\forall x \in \mathbb{R}^{n} \backslash\{0\}, \quad\|x\|_{\mathbf{Y}_{p}^{n}}=\|x\|_{\ell_{p}^{n}} \quad \text { and } \quad \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}} B_{\mathbf{Y}_{p}^{n}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}_{p}^{n}}\right)} \lesssim n^{\frac{1}{p}} . \tag{39}
\end{equation*}
$$

The case $p=\infty$ of Theorem 24 implies Theorem 20through an application of Theorem 21 Indeed, fix $p \geqslant 1$ and $n \in \mathbb{N}$. Suppose that $x, y \in\left(\ell_{p}^{n}\right) \leqslant k$ for some $k \in\{1, \ldots, n\}$. Then $x-y$ has at most $2 k$ nonzero coordinates. Therefore, if $\mathbf{Y}_{\infty}^{n}$ is as in Theorem [24, then by Hölder's inequality we have

$$
\begin{equation*}
(2 k)^{-\max \left\{\frac{1}{2}-\frac{1}{p}, 0\right\}}\|x-y\|_{\ell_{2}^{n}} \leqslant\|x-y\|_{\ell_{p}^{n}} \leqslant(2 k)^{\frac{1}{p}}\|x-y\|_{\ell_{\infty}^{n}}=k^{\frac{1}{p}}\|x-y\|_{\mathbf{Y}_{\infty}^{n}} . \tag{40}
\end{equation*}
$$

Theorem 20 follows by substituting these bounds and the case $p=\infty$ of (39) into (28). Observe that we would have obtained the weaker bound e( $\left.\left(\ell_{p}^{n}\right)_{\leqslant k}\right) \lesssim k^{1 / p+1 / 2}$ if we used (34) instead of (28).

If $p=O(1)$, then one can take $\mathbf{Y}_{p}^{n}=\ell_{p}^{n}$ in Theorem 24. In fact, the direction $z \in S^{n-1}$ at which

$$
\begin{equation*}
\max _{z \in S S^{n-1}} \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z^{\perp}} B \ell_{p}^{n}\right) \tag{41}
\end{equation*}
$$

is attained was determined by Barthe and the author in [BN02]. This result implies that

$$
\begin{equation*}
\forall p \geqslant 1, \quad \max _{z \in S^{n-1}} \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z^{\perp}} B_{\ell_{p}^{n}}\right)}{\operatorname{vol}_{n}\left(B_{\ell_{p}^{n}}\right)}=n^{\frac{1}{p}} \sqrt{\min \{p, n\}} . \tag{42}
\end{equation*}
$$

As (BN02] computes (41) exactly, the implicit constant factors in (42) can be evaluated, but in the present context such precision is of secondary importance. While 42) follows from [BN02] (see the deduction in (Nao17a]), we will give a self-contained proof of (42) in Section 6 as a special case of a more general result that we will use for other purposes as well. In the range $q \in(2, \infty)$, a different approach to computing (41) was found in [KRZ04]. Earlier methods for estimating (41) with worse lower order factors are due to [Sch89] and [Mül90]; the latter is an adaptation of an idea (used for related purposes) in Bou87].

For each $k \in\{1, \ldots, n\}$, by applying (28) with $\mathbf{Y}=\ell_{q}^{n}$ for some $q \geqslant p$, using (42) with $p$ replaced by $q$, and optimizing the resulting bound over $q$, one obtains a result that matches Theorem 20 up to unbounded lower order factors. More precisely, the best that one can get with this approach (up to universal constant factors) is when $q=\max \{2 \log (n / k), p\}$ if $p \leqslant \log (2 k)$. If $p \geqslant \log (2 k)$, then use (28) with $\mathbf{Y}=\ell_{\log (2 k)}^{n}$.

Theorem [24 provides an auxiliary space $\mathbf{Y}$ for which a use of (28) removes the above lower order factors, and yields a sharp result when $p=\infty$ (we conjecture that it is sharp for any $p \geqslant 2$ ). Regardless of whether we apply (28) with the space $\mathbf{Y}=\mathbf{Y}_{\infty}^{n}$ of Theorem 24 or with $\mathbf{Y}=\ell_{q}^{n}$ for a suitable choice of $q \geqslant p$, we have seen that without using an auxiliary space $\mathbf{Y} \neq \ell_{p}^{n}$ in (28) we do not come close to such results.

Even though in Theorem 21 we are interested in extending functions that are Lipschitz in the metric that is induced by the given norm $\|\cdot\|_{\mathbf{X}}$, the underlying reason for the bounds of Theorem 21 is a partitioning scheme (to be described below) that iteratively carves out balls in the metric that is induced by the auxiliary norm $\|\cdot\|_{\mathbf{Y}}$. So, the perturbation of $\mathbf{X}$ into $\mathbf{Y}$ amounts to exhibiting a Lipschitz extension operator through the use of a multi-scale construction that utilizes geometric shapes that differ from balls in the ambient metric. This strategy is feasible because the quantity $e\left(\mathcal{C}_{\mathbf{x}}\right)$ in the left hand side of (32) is a bi-Lipschitz invariant, while the volumes that appear in the right hand side of (32) scale exponentially in $n$. Hence, by passing to an equivalent norm one could hope to reduce the right hand side of (32) significantly, while not changing the left hand side of (32) by too much.

This perturbative approach is decisively useful for $\mathbf{X}=\ell_{\infty}^{n}$. When one unravels the ensuing proofs, the upper bound on $\mathrm{e}\left(\ell_{\infty}^{n}\right)$ of Theorem 14 arises from a multi-scale construction of an extension operator (using a gentle partition of unity [LN05]) that utilizes a partition of space that is obtained by iteratively removing sets of the form $x+r B_{\mathbf{Y}_{\infty}^{n}}$, where $\mathbf{Y}_{\infty}^{n}$ is as in Theorem 24. If one carries out the same procedure while using balls of the intrinsic metric of $\ell_{\infty}^{n}$ (namely, hypercubes $x+r[-1,1]^{n}$ in place of $x+r B_{\mathbf{Y}_{\infty}^{n}}$, which look like hypercubes with "rounded corners"), then only the weaker bound e( $\left.\ell_{\infty}^{n}\right) \lesssim n$ is obtained. We already mentioned that such a phenomenon even occurs in the proof of the Euclidean estimate (27).

The following two examples describe further uses of Theorem 21; we will work out several more later.
Example 25. In the forthcoming work [NS21a], the author and Schechtman prove (for an application to metric embedding theory) the following asymptotic evaluation of the maximal volumes of hyperplane projections of the unit balls of the Schatten-von Neumann trace classes.

$$
\begin{equation*}
\forall q \geqslant 1, \quad \max _{A \in \mathrm{M}_{n}(\mathbb{R}) \backslash\{0\}} \frac{\operatorname{vol}_{n^{2}-1}\left(\operatorname{Proj}_{A^{\perp}} B_{S_{q}^{n}}\right)}{\operatorname{vol}_{n^{2}}\left(B_{S_{q}^{n}}\right)}=n^{\frac{1}{2}+\frac{1}{q}} \sqrt{\min \{q, n\}} . \tag{43}
\end{equation*}
$$

Upon substitution into Theorem 21, this yields the following new estimates on the Lipschitz extension moduli of Schatten-von Neumann trace classes, which holds for every $p \geqslant 1$ and every integer $n \geqslant 2$.

$$
\mathrm{e}\left(\mathrm{~S}_{p}^{n}\right) \lesssim \begin{cases}n^{\frac{1}{2}+\frac{1}{p}} & \text { if } p \in[1,2]  \tag{44}\\ n \sqrt{\min \{p, \log n\}} & \text { if } p \in[2, \infty]\end{cases}
$$

Indeed, by Hölder's inequality $\|\cdot\|_{S_{2}^{n}} \leqslant n^{\max \left\{0, \frac{1}{2}-\frac{1}{p}\right\}_{\|}}\| \|_{S_{p}^{n}}$, so (44) for $p \leqslant \log n$ follows from a substitution of these point-wise bounds and (43) when $q=p$ into the case $\mathbf{X}=\mathbf{Y}=\mathrm{S}_{p}^{n}$ of Theorem 21. The case $p \geqslant \log n$ of (44) follows from the same reasoning using (43) when $q=\log n$ and Theorem 21 for $\mathbf{X}=\mathrm{S}_{p}^{n}$ and $\mathbf{Y}=\mathrm{S}_{q}^{n}$, since in this case $d_{\mathrm{BM}}\left(\mathrm{S}_{p}^{n}, \mathrm{~S}_{q}^{n}\right) \lesssim 1$. Note that, since $\operatorname{dim}\left(\mathrm{S}_{p}^{n}\right)=n^{2}$, for every $p \in[1, \infty]$ the bound on e $\left(\mathrm{S}_{p}^{n}\right)$ in (44] is $o\left(\operatorname{dim}\left(\mathrm{~S}_{p}^{n}\right)\right)$, i.e., it is asymptotically better than what follows from [JLS86].

More generally, given $p \geqslant 1$, an integer $n \geqslant 2$ and $r \in\{3, \ldots, n\}$, let $\left(S_{p}^{n}\right) \leqslant r$ be the set of $n$ by $n$ matrices of rank at most $r$, equipped with the metric inherited from $\mathrm{S}_{p}^{n}$. Then, (44) has the following strengthening.

$$
\mathrm{e}\left(\left(\mathrm{~S}_{p}^{n}\right) \leqslant r\right) \lesssim r^{\max \left\{\frac{1}{p}, \frac{1}{2}\right\}} \sqrt{n} \cdot \begin{cases}\sqrt{\max \left\{\log \left(\frac{n}{r}\right), p\right\}} & \text { if } p \leqslant \log r  \tag{45}\\ \sqrt{\log n} & \text { if } p \geqslant \log r .\end{cases}
$$

To justify (45), apply Theorem 21 with $\mathbf{X}=\mathrm{S}_{p}^{n}$ and $\mathbf{Y}=\mathrm{S}_{q}^{n}$ for some $q \geqslant p$ while using (43), and optimize the resulting bound over $q$. Specifically, since for any $A, B \in\left(S_{p}^{n}\right) \leqslant r$ the matrix $A-B$ has at most $2 r$ nonzero singular values, by Hölder's inequality we have

In combination with (43), we therefore get the following bound from (28).

The $q \geqslant p$ that minimizes the right hand side of (46) is $\max \{2 \log (n / r), p\}$, yielding (45) when $p \leqslant \log r$. If $p \geqslant \log r$, then $\|A-B\|_{S_{p}^{n}}=\|A-B\|_{S_{\log r}^{n}}$ for every $A, B \in\left(S_{p}^{n}\right)_{\leqslant r}$, so (45) reduces to its special case $p=\log r$.

We conjecture that it is possible to replace the logarithmic factor in (45) by a universal constant, i.e.,

$$
\begin{equation*}
\mathrm{e}\left(\left(\mathrm{~S}_{p}^{n}\right) \leqslant r\right) \lesssim r^{\max \left\{\frac{1}{p}, \frac{1}{2}\right\}} \sqrt{n} . \tag{47}
\end{equation*}
$$

As we will see in Section 1.6, Conjecture 26 below is equivalent to the symmetric isomorphic reverse isoperimetry conjecture (see Conjecture 46 for $\mathrm{M}_{n}(\mathbb{R})$ equipped with the operator norm, which is an especially interesting special case of this much more general conjectural phenomenon; by reasoning as we did in the above deduction of Theorem 20 from (the special case $p=\infty$ of) Theorem 24 (recall the discussion immediately following (40), a positive answer to Conjecture 26 would imply (47).

Conjecture 26. For every $n \in \mathbb{N}$ there exists a normed space $\mathbf{Y}=\left(\mathrm{M}_{n}(\mathbb{R}),\|\cdot\|_{\mathbf{Y}}\right)$ such that for every nonzero $n$ by $n$ matrix $A \in \mathrm{M}_{n}(\mathbb{R}) \backslash\{0\}$ we have $\|A\|_{\mathbf{Y}}=\|A\|_{S_{\infty}^{n}}$ and $\operatorname{vol}_{n^{2}-1}\left(\operatorname{Proj}_{A^{\perp}} B_{\mathbf{Y}}\right) \lesssim \operatorname{vol}_{n^{2}}\left(B_{\mathbf{Y}}\right) \sqrt{n}$.

Example 27. Since the $\ell_{\infty}^{n}\left(\ell_{\infty}^{n}\right)$ norm on $\mathrm{M}_{n}(\mathbb{R})$ is isometric to $\ell_{\infty}^{n^{2}}$, by Theorem 24 there is a normed space $\mathbf{Y}=\left(\mathrm{M}_{n}(\mathbb{R}),\|\cdot\|_{\mathbf{Y}}\right)$ that satisfies $\|A\|_{\ell_{\infty}^{n}\left(\ell_{\infty}^{n}\right)} \leqslant\|A\|_{\mathbf{Y}} \lesssim\|A\|_{\ell_{\infty}^{n}\left(\ell_{\infty}^{n}\right)}$ for every $A \in \mathrm{M}_{n}(\mathbb{R})$, and

$$
\max _{A \in \mathrm{M}_{n}(\mathbb{R}) \backslash\{0\}} \frac{\operatorname{vol}_{n^{2}-1}\left(\operatorname{Proj}_{A^{\perp}} B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n^{2}}\left(B_{\mathbf{Y}}\right)}=O(1) .
$$

By Hölder's inequality, for every $p, q \in[1, \infty]$ and $A \in \mathrm{M}_{n}(\mathbb{R})$ we have

$$
\left.\|A\|_{\ell_{p}^{n}\left(\ell_{q}^{n}\right)} \leqslant n^{\frac{1}{p}+\frac{1}{q}}\|A\|_{\ell_{\infty}^{n}\left(\ell_{\infty}^{n}\right)} \leqslant n^{\frac{1}{p}+\frac{1}{q}}\|A\|_{\mathbf{Y}} \quad \text { and } \quad\|A\|_{\ell_{2}^{n}\left(\ell_{2}^{n}\right)} \leqslant n^{\max \left\{\frac{1}{2}-\frac{1}{p}, 0\right\}+\max \left\{\frac{1}{2}-\frac{1}{q}, 0\right\}}\|A\|_{\rho}^{n} \ell_{q}^{n}\right) .
$$

Therefore, Theorem 21 gives the Lipschitz extension bound

$$
\begin{equation*}
\mathrm{e}\left(\ell_{p}^{n}\left(\ell_{q}^{n}\right)\right) \lesssim n^{\frac{1}{p}+\frac{1}{q}+\max \left\{\frac{1}{2}-\frac{1}{p}, 0\right\}+\max \left\{\frac{1}{2}-\frac{1}{q}, 0\right\}}=n^{\max \left\{1, \frac{1}{p}+\frac{1}{q}, \frac{1}{2}+\frac{1}{p}, \frac{1}{2}+\frac{1}{q}\right\}} . \tag{48}
\end{equation*}
$$

As in the case of $\ell_{p}^{n}$, we get (48) if $p, q=O(1)$ by using Theorem 21 with $\mathbf{Y}=\mathbf{X}=\ell_{p}^{n}\left(\ell_{p}^{n}\right)$, but otherwise we need to work with an auxiliary space $\mathbf{Y} \neq \mathbf{X}$ as above. Specifically, in Section 6 we will prove the following
asymptotic evaluation of the maximal volume of hyperplane projections of the unit ball of $\ell_{p}^{n}\left(\ell_{q}^{n}\right)$ :

$$
\max _{A \in \mathrm{M}_{n}(\mathbb{R}) \backslash\{0\}} \frac{\operatorname{vol}_{n^{2}-1}\left(\operatorname{Proj}_{A^{\perp}} B_{\ell_{p}^{n}\left(\ell_{q}^{n}\right)}\right)}{\operatorname{vol}_{n^{2}}\left(B_{\ell_{p}^{n}\left(\ell_{q}^{n}\right)}\right)}= \begin{cases}n & \text { if } n \leqslant \min \{\sqrt{p}, q\},  \tag{49}\\ \sqrt{q} n^{\frac{1}{2}+\frac{1}{q}} & \text { if } q \leqslant n \leqslant \sqrt{p}, \\ \sqrt{p} & \text { if } \sqrt{p} \leqslant n \leqslant \min \{p, q\}, \\ \sqrt{p q} n^{\frac{1}{q}-\frac{1}{2}} & \text { if } \max \{\sqrt{p}, q\} \leqslant n \leqslant p, \\ n^{\frac{1}{2}+\frac{1}{p}} & \text { if } p \leqslant n \leqslant q, \\ \sqrt{q} n^{\frac{1}{p}+\frac{1}{q}} & \text { if } n \geqslant \max \{p, q\} .\end{cases}
$$

The intricacy of (49) is perhaps unexpected, though it is nonetheless sharp in all of the six ranges (depending on the relative locations of $p, q, n$ and, somewhat curiously, $\sqrt{p}$ ) that appear in (49). By reasoning analogously to the discussion following (42), one can prove a bound on $\mathrm{e}\left(\ell_{p}^{n}\left(\ell_{q}^{n}\right)\right)$ that matches (48) up to lower order factors by applying Theorem 21 with $\mathbf{Y}=\ell_{r}^{n}\left(\ell_{s}^{n}\right)$ and then optimizing over $r, s \geqslant 1$. For the sole purpose of this application, only the range $n \geqslant \max \{p, q\}$ of (49) is needed. However, results such as (49) have geometric interest in their own right for all of the possible values of the relevant parameters. We will actually prove a version of (49) for $\ell_{p}^{n}\left(\ell_{q}^{m}\right)$ even when $n \neq m$; the case of rectangular matrices is independently interesting, but we will also use it elsewhere (see Remark[55bbelow).

Problem 28. Determine the exact maximizers of volumes of hyperplane projections of the unit balls of $S_{p}^{n}$ and $\ell_{p}^{n}\left(\ell_{q}^{n}\right)$, i.e., for which $A \in \mathrm{M}_{n}(\mathbb{R}) \backslash\{0\}$ are the maxima in (43) and (49) attained.
1.5. A dimension-independent extension theorem. In the preceding sections we stated all of the extension theorems using the traditional setup that aims to extend a Lipschitz function to a function that is Lipschitz with respect to the given metric. However, all of our new (positive) extension theorems are a consequence of Theorem 29] below, which is a nonstandard Lipschitz extension theorem.

Theorem 29asserts that if $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a normed space and $f$ is a 1-Lipschitz function from a subset of $\mathbb{R}^{n}$ to a Banach space $\mathbf{Z}$, then $f$ can be extended to a $\mathbf{Z}$-valued function that is defined on all of $\mathbb{R}^{n}$ and is $O(1)$-Lipschitz with respect to the metric that is induced on $\mathbb{R}^{n}$ by the norm $\|\|\cdot\|\|=2\|\cdot\|_{\Pi^{*} \mathbf{X}} / \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)$, i.e., a suitable rescaling of the norm whose unit ball is the polar projection body of $B_{\mathbf{X}}$. This rescaling ensures that $\|\|\cdot\|\|$ dominates $\|\cdot\|_{\mathbf{x}}$; indeed, by an elementary geometric argument (see Remark 111),

$$
\begin{equation*}
\forall x \in \mathbb{R}^{n}, \quad\|x\|_{\mathbf{X}} \leqslant \frac{2\|x\|_{\Pi^{*} \mathbf{X}}}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} \leqslant n\|x\|_{\mathbf{X}} . \tag{50}
\end{equation*}
$$

Thus, the conclusion of Theorem 29 that the extended function is Lipschitz with respect to \|II•\|| is less stringent than the traditional requirement that it should be Lipschitz with respect to $\|\cdot\|_{\mathbf{X}}$, but Theorem 29 has the feature that the upper bound on the Lipschitz constant is independent of the dimension.

Theorem 29. Fix $n \in \mathbb{N}$, a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{x}}\right)$ and a Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$. Suppose that $\mathcal{C} \subseteq \mathbb{R}^{n}$ and $f: \mathcal{C} \rightarrow \mathbf{Z}$ is 1 -Lipschitz with respect to the metric that is induced by $\|\cdot\|_{\mathbf{X}}$, i.e., $\|f(x)-f(y)\|_{\mathbf{Z}} \leqslant\|x-y\|_{\mathbf{X}}$ for every $x, y \in \mathcal{C}$. Then, there exists $F: \mathbb{R}^{n} \rightarrow \mathbf{Z}$ that coincides with $f$ on $\mathcal{C}$ and satisfies

$$
\forall x, y \in \mathbb{R}^{n}, \quad\|F(x)-F(y)\|_{\mathbf{Z}} \lesssim \frac{\|x-y\|_{n^{*} \mathbf{x}}}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} .
$$

To see how Theorem 29implies Theorem 21, denote (in the setting of the statement of Theorem 21):

$$
\begin{equation*}
M=\sup _{\substack{x, y \in \mathcal{e} \\ x \neq y}}\left(\frac{\|x-y\|_{\mathbf{X}}}{\|x-y\|_{\mathbf{Y}}}\right) \quad \text { and } \quad M^{\prime}=\sup _{\substack{x, y \in \mathcal{C} \\ x \neq y}}\left(\frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{(x-y) \perp} B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \cdot \frac{\|x-y\|_{\ell_{2}^{n}}}{\|x-y\|_{\mathbf{X}}}\right) . \tag{51}
\end{equation*}
$$

Thus, every $x, y \in \mathcal{C}$ satisfy $\|x-y\|_{\mathbf{X}} \leqslant M\|x-y\|_{\mathbf{Y}}$ and, recalling (30), also $\|x-y\|_{\Pi^{*} \mathbf{Y}} / \operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right) \leqslant M^{\prime}\|x-y\|_{\mathbf{X}}$. Let $\left(\mathbf{Z},\|\cdot\|_{\mathbf{z}}\right)$ be any Banach space and consider an arbitrary subset $\mathcal{C}^{\prime} \subseteq \mathcal{C}$. If $f: \mathcal{C}^{\prime} \rightarrow \mathbf{Z}$ is 1 -Lipschitz with respect to the metric that is induced by $\|\cdot\|_{\mathbf{X}}$, then $f / M$ is l-Lipschitz with respect to the metric that is induced by $\mathbf{Y}$. By Theorem 29 (with $\mathbf{X}$ replaced by $\mathbf{Y}, \mathcal{C}$ replaced by $\mathfrak{C}^{\prime}, f$ replaced by $f / M$ ) we therefore
see that there exists $F: \mathbb{R}^{n} \rightarrow \mathbf{Z}$ (for Theorem 21 we only need $F$ to be defined on $\mathcal{C}$ ) that extends $F$ and satisfies $\|F(x)-F(y)\|_{\mathbf{Z}} \lesssim M\|x-y\|_{\Pi^{*} \mathbf{Y}} / \operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right) \leqslant M M^{\prime}\|x-y\|_{\mathbf{X}}$ for all $x, y \in \mathcal{C}$. This coincides with (28).

Remark 30. Given $p \geqslant 1$, consider what happens when we apply Theorem 29 to the space $\mathbf{Y}_{p}^{n}$ of Theorem24. We get that for any $\mathcal{C} \subseteq \mathbb{R}^{n}$ and any Banach space $\mathbf{Z}$, if $f: \mathcal{C} \rightarrow \mathbf{Z}$ is 1 -Lipschitz with respect to the $\ell_{p}^{n}$ metric, then $f$ can be extended to $F: \mathbb{R}^{n} \rightarrow \mathbf{Z}$ that is $O\left(n^{1 / p}\right)$-Lipschitz with respect to the Euclidean metric. When $p<2$, the Lipschitz assumption on $f$ is less stringent than requiring it to be $O(1)$-Lipschitz with respect to the Euclidean metric, but we then get an extension $F$ that is $O\left(n^{1 / p}\right)$-Lipschitz with respect to the Euclidean metric; this upper bound on the Lipschitz constant of $F$ is asymptotically larger than the $O(\sqrt{n})$ bound that we would get if $f$ were assumed to be 1-Lipschitz with respect to the Euclidean metric and we applied the second inequality in (22), but we get it while requiring less from $f$. In particular, when $p=1$ we see that any $\mathbf{Z}$-valued function on a subset of $\mathbb{R}^{n}$ that is 1-Lipschitz with respect to the $\ell_{1}^{n}$ metric can be extended to a $\mathbf{Z}$-valued function defined on all of $\mathbb{R}^{n}$ whose Lipschitz constant with respect to the Euclidean metric is $O(n)$, while an application of [JLS86] will give an extension that is $O(n)$-Lipschitz with respect to the $\ell_{1}^{n}$ metric. On the other hand, if $p>2$, then the Lipschitz assumption on $f$ is more stringent than requiring it to be $O(1)$-Lipschitz with respect to the Euclidean metric, but we then get an extension $F$ that is $O\left(n^{1 / p}\right)$-Lipschitz with respect to the Euclidean metric, which is asymptotically better than the $O(\sqrt{n})$ bound from (22). In particular, when $p=\infty$ we see that any $\mathbf{Z}$-valued function on a subset of $\mathbb{R}^{n}$ that is 1-Lipschitz with respect to the $\ell_{\infty}^{n}$ metric can be extended to a $\mathbf{Z}$-valued function on all of $\mathbb{R}^{n}$ whose Lipschitz constant with respect to the Euclidean metric is $O(1)$.
1.6. Isomorphic reverse isoperimetry. All of the applications that we found for Theorem 21 proceed by bounding the volumes of hyperplane projections of $B_{\mathbf{Y}}$ that appear in right hand side of (28) by

$$
\begin{equation*}
\operatorname{MaxProj}\left(B_{\mathbf{Y}}\right) \stackrel{\text { def }}{=} \max _{z \in S^{n-1}} \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z^{\perp}} B_{\mathbf{Y}}\right) \tag{52}
\end{equation*}
$$

Thus, it follows from (29) that for any two normed spaces $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right), \mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ with $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$,

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim \frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right) \tag{53}
\end{equation*}
$$

While there could conceivably be an application of (29) that is more refined than (53), in this section we will investigate the ramifications of bounding $\operatorname{MaxProj}\left(B_{\mathbf{X}}\right)$ as a way to use Theorem 21. This will relate to the isomorphic reverse isoperimetric phenomena that we conjectured in Section 1.1.1.

Any origin-symmetric convex body $L \subseteq \mathbb{R}^{n}$ satisfies

$$
\begin{equation*}
\operatorname{MaxProj}(L) \gtrsim \frac{\operatorname{vol}_{n-1}(\partial L)}{\sqrt{n}} \tag{54}
\end{equation*}
$$

Indeed, this follows immediately from the following classical Cauchy surface area formula (see e.g. equation 5.73 in Sch14]) by bounding the integrand by its maximum.

$$
\begin{equation*}
\operatorname{vol}_{n-1}(\partial L)=\frac{2 \sqrt{\pi} \Gamma\left(\frac{n+1}{2}\right)}{\Gamma\left(\frac{n}{2}\right)} f_{S^{n-1}} \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z^{\perp}} L\right) \mathrm{d} z=\sqrt{n} f_{S^{n-1}} \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z^{\perp}} L\right) \mathrm{d} z \tag{55}
\end{equation*}
$$

Remark 31. Using (54), Theorem 24 implies that Conjecture 9 (isomorphic reverse isoperimetry) holds (with $S$ the identity mapping) when $K=B_{\ell_{p}^{n}}$ for any $p \geqslant 1$ and $n \in \mathbb{N}$. Indeed, let $\mathbf{Y}_{p}^{n}$ be the normed space from Theorem 24. By the first inequality in (40) we have

$$
\begin{equation*}
\operatorname{vol}_{n}\left(B_{\mathbf{Y}_{p}^{n}}\right)^{\frac{1}{n}}=\operatorname{vol}_{n}\left(B_{\ell_{p}^{n}}\right)^{\frac{1}{n}}=n^{-\frac{1}{p}}, \tag{56}
\end{equation*}
$$

where the last equivalence in (56) is a standard computation (e.g. [Pis89, page 11]). By (54) and (56), the second inequality in (40) implies that the isoperimetric quotient of $B_{\mathbf{Y}_{p}^{n}}$ is $O(\sqrt{n})$. So, Conjecture 9 holds for $K=B_{\ell_{p}^{n}}$ if we take $L$ to be a rescaling by a universal constant factor of $B_{\mathbf{Y}_{p}^{n}}$ so that $L \subseteq K$.

Thanks to (54), if we set $K=B_{\mathbf{X}}$ and $L=B_{\mathbf{Y}}$ in (53), then the right hand side of (53) satisfies

$$
\begin{equation*}
\frac{\operatorname{MaxProj}(L)}{\operatorname{vol}_{n}(L)} \operatorname{diam}_{\ell_{2}^{n}}(K) \gtrsim \frac{\operatorname{vol}_{n-1}(\partial L)}{\sqrt{n} \operatorname{vol}_{n}(L)} \operatorname{diam}_{\ell_{2}^{n}}(K)=\frac{\operatorname{iq}(L)}{\sqrt{n}} \cdot \frac{\operatorname{diam}_{\ell_{2}^{n}(K)}}{\operatorname{vol}_{n}(L)^{\frac{1}{n}}} \gtrsim \frac{\operatorname{diam}_{\ell_{2}^{n}(K)}}{\operatorname{vol}_{n}(K)^{\frac{1}{n}}}, \tag{57}
\end{equation*}
$$

where we recall notation (11) for the isoperimetric quotient $\mathrm{iq}(\cdot)$ and the last step uses the isoperimetric theorem (12) and the assumption $L \subseteq K$. The following proposition explains what it would entail for one to be able to reverse (57) after an application of a suitable linear transformation; in particular, it shows that one can find $S \in S L_{n}(\mathbb{R})$ and an origin-symmetric convex body $L \subseteq S K$ such that

$$
\frac{\operatorname{MaxProj}(L)}{\operatorname{vol}_{n}(L)} \operatorname{diam}_{\ell_{2}^{n}}(S K) \lesssim \frac{\operatorname{diam}_{\ell_{2}^{n}}(S K)}{\operatorname{vol}_{n}(K)^{\frac{1}{n}}}
$$

if and only if Conjecture 10 on weak isomorphic reverse isoperimetry holds for $K$.
Proposition 32. The following two statements are equivalent for every $n \in \mathbb{N}$, every origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$ and every $\alpha>0$.
(1) There exist a linear transformation $S \in S_{n}(\mathbb{R})$ and an origin-symmetric convex body $L \subseteq S K$ with

$$
\begin{equation*}
\frac{\operatorname{MaxProj}(L)}{\operatorname{vol}_{n}(L)} \operatorname{vol}_{n}(K)^{\frac{1}{n}} \lesssim \alpha . \tag{58}
\end{equation*}
$$

(2) There exist a linear transformation $S \in S L_{n}(\mathbb{R})$ and an origin-symmetric convex body $L \subseteq S K$ that satisfies $\sqrt[n]{\operatorname{vol}_{n}(L)} \geqslant \beta \sqrt[n]{\operatorname{vol}_{n}(K)}$ and $\mathrm{iq}(L) \leqslant \gamma \sqrt{n}$ for some $\beta \gtrsim 1 / \alpha$ and $\gamma \lesssim \alpha$ with $\gamma / \beta \lesssim \alpha$.
Proof. For the implication (1) $\Longrightarrow(2)$ denote $\beta=\sqrt[n]{\operatorname{vol}_{n}(L)} / \sqrt[n]{\operatorname{vol}_{n}(K)}$ and $\gamma=\mathrm{iq}(L) / \sqrt{n}$. Then

$$
\alpha \stackrel{58}{\gtrsim} \frac{\operatorname{MaxProj}(L)}{\operatorname{vol}_{n}(L)} \operatorname{vol}_{n}(K)^{\frac{1}{n}} \stackrel{\text { 54] }}{\gtrsim} \frac{\operatorname{vol}_{n-1}(\partial L)}{\operatorname{vol}_{n}(L) \sqrt{n}} \operatorname{vol}_{n}(K)^{\frac{1}{n}}=\frac{\gamma}{\beta} .
$$

Since by the isoperimetric theorem (12) we have $\gamma \gtrsim 1$, it follows from this that $\beta \gtrsim 1 / \alpha$, and since $L \subseteq S K$ and $S \in S L_{n}(\mathbb{R})$, we have $\operatorname{vol}_{n}(L) \leqslant \operatorname{vol}_{n}(K)$, so $\beta \leqslant 1$ and it also follows from this that $\gamma \lesssim \alpha$.

For the implication (2) $\Longrightarrow(1)$, fix $T \in \mathrm{SL}_{n}(\mathbb{R})$ with $\operatorname{vol}_{n-1}(\partial T L)=\min \left\{\operatorname{vol}_{n-1}\left(\partial T^{\prime} L\right): T^{\prime} \in \mathrm{SL}_{n}(\mathbb{R})\right\}$, i.e., $T L$ is in its minimum surface area position [Pet61]. By definition, $\operatorname{vol}_{n-1}(\partial T L) \leqslant \operatorname{vol}_{n-1}(\partial L)$ and by Proposition 3.1 in the work GP99] of Giannopoulos and Papadimitrakis combined with (54) we have

$$
\operatorname{MaxProj}(T L)=\frac{\operatorname{vol}_{n-1}(\partial T L)}{\sqrt{n}}
$$

Consequently, if $L$ satisfies part (2) of Proposition 32, then

$$
\frac{\operatorname{MaxProj}(T L)}{\operatorname{vol}_{n}(T L)} \operatorname{vol}_{n}(K)^{\frac{1}{n}}=\frac{\operatorname{vol}_{n-1}(\partial T L)}{\operatorname{vol}_{n}(T L) \sqrt{n}} \operatorname{vol}_{n}(K)^{\frac{1}{n}} \leqslant \frac{\operatorname{vol}_{n-1}(\partial L)}{\operatorname{vol}_{n}(T L) \sqrt{n}} \operatorname{vol}_{n}(K)^{\frac{1}{n}}=\frac{\operatorname{iq}(L)}{\sqrt{n}}\left(\frac{\operatorname{vol}_{n}(K)}{\operatorname{vol}_{n}(L)}\right)^{\frac{1}{n}} \leqslant \frac{\gamma}{\beta} \lesssim \alpha
$$

It follows that (1) holds with $S$ replaced by $T S \in \mathrm{SL}_{n}(\mathbb{R})$ and $L$ replaced by $T L \subseteq T S K$.
Since when $\alpha \lesssim 1$ in Proposition 32 the assertion of its part (2) coincides with Conjecture 10 it follows that Conjecture 10, and a fortiori Conjecture 9 , imply that for any normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ there is $S \in \mathrm{SL}_{n}(\mathbb{R})$ such that $\mathrm{e}(\mathbf{X})$ is at most a universal constant multiple of $\operatorname{diam}_{\ell_{2}^{n}}\left(S B_{\mathbf{X}}\right) / \sqrt[n]{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}$. Indeed, this follows by applying Theorem 21 to the normed spaces $\mathbf{X}^{\prime}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}^{\prime}}\right)$ and $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ whose unit balls are $S B_{\mathbf{X}}$ and $L$, respectively, where $S$ and $L$ are as in part (1) of Proposition 32 for $K=B_{\mathbf{X}}$, while noting that $\mathrm{e}\left(\mathbf{X}^{\prime}\right)=e(\mathbf{X})$ since $\mathbf{X}^{\prime}$ is isometric to $\mathbf{X}$. We record this conclusion as the following corollary.

Corollary 33. If Conjecture 10 holds for a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$, then there is $S \in S L_{n}(\mathbb{R})$ such that

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim \frac{\operatorname{diam}_{\ell_{2}^{n}}\left(S B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{\frac{1}{n}}} \tag{59}
\end{equation*}
$$

The upshot of Corollary 33 is that the right hand side of (59) involves only Euclidean diameters and $n$ 'th roots of volumes, which are typically much easier to estimate than extremal volumes of hyperplane projections. This comes at the cost of having to find the auxiliary linear transformation $S \in S L_{n}(\mathbb{R})$, but we expect that in concrete settings it will be simple to determine $S$. Moreover, in all of the specific examples of spaces for which we are interested (at least initially) in estimating their Lipschitz extension modulus, $S$ should be the identity mapping. We will discuss this matter and its consequences in Section 1.6.2.

Remark 34. There is a degree of freedom that the above discussion does not exploit. Let $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ be a normed space. By (31), we know that $e(\mathbf{X})$ is bounded from above by a constant multiple of the minimum of $\operatorname{diam}_{\Pi^{*} \mathbf{Y}}\left(B_{\mathbf{X}}\right) / \operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)$ over all the normed spaces $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ for which $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$. By (53), to control this minimum it suffices to estimate the minimum of $\operatorname{MaxProj}\left(\mathrm{B}_{\mathbf{Y}}\right) / \operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)$ over all such $\mathbf{Y}$, which relates to isomorphic reverse isoperimetric phenomena. But, we could also take a normed space $\mathbf{W}=\left(\mathbb{R}^{m},\|\cdot\|_{\mathbf{W}}\right)$ for $m \geqslant n$ such that $B_{\mathbf{W}} \cap \mathbb{R}^{n}=B_{\mathbf{X}}$ (we need that $\mathbf{W}$ contains an isometric copy of $\mathbf{X}$ ), estimate either of the two minima above for the super-space $\mathbf{W}$, and then use $e(\mathbf{X}) \leqslant e(\mathbf{W})$. Thus, it would suffice to embed $\mathbf{X}$ into a larger normed space that exhibits good isomorphic reverse isoperimetry. Our conjectures imply that such an embedding step is not needed, namely we expect that the desired isomorphic reverse isoperimetric property holds for $\mathbf{X}$. Nevertheless, it could be that by finding a suitable super-space $\mathbf{W}$ one could bound $\mathrm{e}(\mathbf{X})$ while circumventing the difficulty of proving Conjecture 10 for $\mathbf{X}$. For example, if $\mathbf{X}$ is a subspace of $\ell_{\infty}^{m}$ for some $m=O(n)$, then by Theorem 14 we know that $\mathrm{e}(\mathbf{X}) \lesssim \sqrt{n}$, but this is because we know that $\ell_{\infty}^{m}$ has the desired isomorphic reverse isoperimetric property, and it is not clear how to prove it for $\mathbf{X}$ itself. It is also unclear how to construct for a given normed $\mathbf{X}$ a super-space $\mathbf{W}$ that could be used as above. We leave the exploration of this possibility for future research.
1.6.1. A spectral interpretation, reverse Faber-Krahn and the Cheeger space of a normed space. We will henceforth quantify the extent to which Conjecture 10 holds through the following condition:

$$
\begin{equation*}
\frac{\mathrm{iq}(L)}{\sqrt{n}}\left(\frac{\operatorname{vol}_{n}(K)}{\operatorname{vol}_{n}(L)}\right)^{\frac{1}{n}}=\frac{\operatorname{vol}_{n}(K)^{\frac{1}{n}}}{\sqrt{n}}\left(\frac{\operatorname{vol}_{n-1}(\partial L)}{\operatorname{vol}_{n}(L)}\right) \leqslant \alpha . \tag{60}
\end{equation*}
$$

The factors $\mathrm{iq}(L) / \sqrt{n}$ and $\left(\operatorname{vol}_{n}(K) / \operatorname{vol}_{n}(L)\right)^{1 / n}$ that appear in the left hand side of (60) are at least a positive universal constant (by, respectively, the isoperimetric theorem and the assumed inclusion $L \subseteq K$ ), so (60) implies that $\sqrt[n]{\operatorname{vol}_{n}(L)} \gtrsim \alpha^{-1} \sqrt[n]{\operatorname{vol}_{n}(K)}$ and iq $(L) \leqslant \alpha \sqrt{n}$. Thus, if $\alpha=O(1)$, then (60) is equivalent to the conclusion of Conjecture 10 . However, even though Conjecture 10 expresses our expectation that (60) is always achievable with $\alpha=O(1)$ upon a judicious choice of the Euclidean structure on $\mathbb{R}^{n}$, in lieu of Conjecture 10 it would still be valuable to obtain (60) with $\alpha$ unbounded but slowly growing. In such a situation, the bi-parameter quantification that we used in part (2) of Proposition 32 contains more geometric information than (60), but below we will work with (60) in order to simplify the ensuing discussion; this suffices for our purposes because (60) is what shows up in all of the applications herein (per the proof Proposition 32) since they all proceed by bounding the right hand side of (53) from above.

Alter and Caselles proved (AC09] that for every convex body $K \subseteq \mathbb{R}^{n}$ there is a unique measurable set $A \subseteq K$, which we call the Cheeger body of $K$ and denote $\operatorname{Ch} K$, satisfying $\operatorname{Per}(A) / \operatorname{vol}_{n}(A) \leqslant \operatorname{Per}(B) / \operatorname{vol}_{n}(B)$ for every measurable $B \subseteq K$ with $\operatorname{vol}_{n}(B)>0$, where $\operatorname{Per}(\cdot)$ denotes perimeter in the sense of Caccioppoli and de Giorgi; this notion is covered in [AFP00] but we do not need to recall its definition here since the perimeter of a convex body coincides with the ( $n-1$ )-dimensional Hausdorff measure of its boundary. It was proved in AC09 that $\mathrm{Ch} K$ is convex and its boundary is $C^{1,1}$. Further information on this remarkable theorem can be found in [AC09], where $\mathrm{Ch} K$ is characterized in terms of the mean curvature of its boundary through the work [ACC05] of Alter, Caselles and Chambolle (see also the precursor [CCN07] by Caselles, Chambolle and Novaga which obtained these statements under stronger assumptions on $K$ ).

Beyond the fact that it allows us to use the notation $\mathrm{Ch} K$ and call it the Cheeger body of $K$, the aforementioned uniqueness statement will be used substantially in the ensuing reasoning. It implies in particular that if $K$ is origin-symmetric, then so is $\operatorname{Ch} K$. Consequently, if $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a normed space, then $\mathrm{Ch} B_{\mathbf{X}}$ is the unit ball of a normed space that we denote by $\mathrm{Ch} \mathbf{X}$ and call the Cheeger space of $\mathbf{X}$.

For a convex body $K \subseteq \mathbb{R}^{n}$, let $\lambda(K)$ be the smallest Dirichlet eigenvalue of the Laplacian on $K$, namely it is the smallest $\lambda>0$ for which there is a nonzero function $\varphi: K \rightarrow \mathbb{R}$ that is smooth on the interior of $K$, vanishes on the boundary of $K$, and satisfies $\Delta \varphi=-\lambda \varphi$ on the interior of $K$; see e.g. [PS51, CH53, Cha84] for background on this classical topic. If $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a normed space, then we denote $\lambda(\mathbf{X})=\lambda\left(B_{\mathbf{X}}\right)$.

The quantity $h(K)=\operatorname{vol}_{n-1}(\partial \operatorname{Ch} K) / \operatorname{vol}_{n}(\operatorname{Ch} K)$ is called the Cheeger constant of $K$; it relates to $\lambda(K)$ by

$$
\begin{equation*}
\frac{2}{\pi} \sqrt{\lambda(K)} \leqslant h(K)=\frac{\operatorname{vol}_{n-1}(\partial \operatorname{Ch} K)}{\operatorname{vol}_{n}(\operatorname{Ch} K)} \leqslant 2 \sqrt{\lambda(K)} . \tag{61}
\end{equation*}
$$

It is important for our purposes that the constants appearing in (61) are independent of the dimension $n$. The second inequality in (61) is the Cheeger inequality for the Dirichlet Laplacian on Euclidean domains. Cheeger's proof of it for compact Riemannian manifolds without boundary appears in [Che70] and that proof works mutatis mutandis in the present setting; see its derivation in e.g. the appendix of [LW97]. The first inequality in (61) can be called the Buser inequality for the Dirichlet Laplacian on convex Euclidean domains, since Buser proved [Bus82] its analogue for compact Riemannian manifolds without boundary that have a lower bound on their Ricci curvature. In our setting, this reverse Cheeger inequality is more recent, namely it was noted for planar convex sets by Parini (Par17] and in any dimension by Brasco [Bra20]. It can be justified quickly using the convexity of $K$ and its Cheeger body $\mathrm{Ch} K$ as follows. By a classical theorem of Pólya we have $\lambda(K) \leqslant \pi^{2}\left(\operatorname{vol}_{n-1}(\partial K) / \mathrm{vol}_{n}(K)\right)^{2} / 4$ (Pólya proved this for planar convex sets, but in [JS82b] Joó and Stachó carried out Pólya's approach for convex bodies in $\mathbb{R}^{n}$ for any $n \in \mathbb{N})$. Therefore, $\lambda(K) \leqslant \lambda(\operatorname{Ch} K) \leqslant \pi^{2}\left(\operatorname{vol}_{n-1}(\partial \operatorname{Ch} K) / \operatorname{vol}_{n}(\operatorname{Ch} K)\right)^{2} / 4=\pi^{2} h(K)^{2} / 4$, since $\operatorname{Ch} K$ is convex.

Let $j_{n / 2-1,1}$ be the smallest positive zero of the Bessel function $J_{n / 2-1}$; see Chapter 4 of [AAR99] for a treatment of Bessel functions and their zeros, though here we will only need to know that $j_{n / 2-1,1}=n$ (see [Tri49] for more precise asymptotics). By classical computations (see e.g. equation 1.29 in [Hen06),

$$
\lambda\left(B_{\ell_{2}^{n}}\right)=j_{\frac{n}{2}-1,1}^{2} .
$$

The Faber-Krahn inequality [Fab23, Kra26] (see also e.g. [PS51, Cha84]) asserts that $\lambda(K)$ is at least the first Dirichlet eigenvalue of a Euclidean ball whose volume is the same as the volume of $K$. Thus,

$$
\lambda(K) \operatorname{vol}_{n}(K)^{\frac{2}{n}} \geqslant \lambda\left(B_{\ell_{2}^{n}}\right) \operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right)^{\frac{2}{n}}=j_{\frac{n}{2}-1,1}^{2} \operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right)^{\frac{2}{n}}=n,
$$

where we used the straightforward fact that $\lambda(r K)=\lambda(K) / r^{2}$ for every $r>0$.
Observe that (61) can be rewritten as follows for every convex body $K \subseteq \mathbb{R}^{n}$.

$$
\frac{2}{\pi}\left(\frac{\lambda(K) \operatorname{vol}_{n}(K)^{\frac{2}{n}}}{n}\right)^{\frac{1}{2}} \leqslant \frac{\mathrm{iq}(\operatorname{Ch} K)}{\sqrt{n}}\left(\frac{\operatorname{vol}_{n}(K)}{\operatorname{vol}_{n}(\operatorname{Ch} K)}\right)^{\frac{1}{n}} \leqslant 2\left(\frac{\lambda(K) \operatorname{vol}_{n}(K)^{\frac{2}{n}}}{n}\right)^{\frac{1}{2}} .
$$

Hence, for every $\alpha>0$ we have

$$
\begin{equation*}
\frac{\mathrm{iq}(\operatorname{Ch} K)}{\sqrt{n}}\left(\frac{\operatorname{vol}_{n}(K)}{\operatorname{vol}_{n}(\operatorname{Ch} K)}\right)^{\frac{1}{n}} \lesssim \alpha \Longleftrightarrow \lambda(K) \operatorname{vol}_{n}(K)^{\frac{2}{n}} \lesssim \alpha^{2} n \tag{62}
\end{equation*}
$$

Since $\mathrm{Ch} K$ is convex, the convex body $L \subseteq K$ that minimizes the left hand side of (60) is equal to $\mathrm{Ch} K$. We therefore see that Conjecture 35 below is equivalent to Conjecture 10. Furthermore, if one of these two conjectures hold for a matrix $S \in S L_{n}(\mathbb{R})$, then the same matrix would work for the other conjecture.

Conjecture 35 (reverse Faber-Krahn). For any origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$ there exists a volumepreserving linear transformation $S \in \mathrm{SL}_{n}(\mathbb{R})$ such that

$$
\lambda(S K) \operatorname{vol}(K)^{\frac{2}{n}}=n
$$

This spectral interpretation of Conjecture 10 is useful for multiple purposes, including the following lemma whose proof appears in Section 6.1. For its statement, as well as throughout the ensuing discussion, recall that a basis $x_{1}, \ldots, x_{n}$ of an $n$-dimensional normed space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ is a 1-unconditional basis of $\mathbf{X}$ if $\left\|\varepsilon_{1} a_{1} x_{1}+\ldots+\varepsilon_{n} a_{n} x_{n}\right\|_{\mathbf{X}}=\left\|a_{1} x_{1}+\ldots+a_{n} x_{n}\right\|_{\mathbf{X}}$ for every choice of scalars $a_{1}, \ldots, a_{n} \in \mathbb{R}$ and signs $\varepsilon_{1}, \ldots, \varepsilon_{n} \in\{-1,1\}$. When we say that $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is an unconditional normed space, we mean that the standard (coordinate) basis $e_{1}, \ldots, e_{n}$ of $\mathbb{R}^{n}$ is a 1-unconditional basis of $\mathbf{X}$.

Lemma 36 (closure of Conjecture 10 under unconditional composition). Fix $n \in \mathbb{N}$ and $m_{1}, \ldots, m_{n} \in \mathbb{N}$. Let $\mathbf{X}_{1}=\left(\mathbb{R}^{m_{1}},\|\cdot\|_{\mathbf{X}_{1}}\right), \ldots, \mathbf{X}_{n}=\left(\mathbb{R}^{m_{n}},\|\cdot\|_{\mathbf{X}_{n}}\right)$ be normed spaces. Also, let $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ be an unconditional normed space. Define a normed space $\mathbf{X}=\left(\mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}},\|\cdot\|_{\mathbf{X}}\right)$ by

$$
\begin{equation*}
\forall x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}}, \quad\|x\|_{\mathbf{X}} \stackrel{\text { def }}{=}\left\|\left(\left\|x_{1}\right\|_{\mathbf{X}_{1}}, \ldots,\left\|x_{n}\right\|_{\mathbf{X}_{n}}\right)\right\|_{\mathbf{E}} \tag{63}
\end{equation*}
$$

Suppose that there exist $\alpha>0$, linear transformations $S_{1} \in \operatorname{SL}_{m_{1}}(\mathbb{R}), \ldots, S_{n} \in S_{m_{n}}(\mathbb{R})$, and normed spaces $\mathbf{Y}_{1}=\left(\mathbb{R}^{m_{1}},\|\cdot\|_{\mathbf{Y}_{1}}\right), \ldots, \mathbf{Y}_{n}=\left(\mathbb{R}^{m_{n}},\|\cdot\|_{\mathbf{Y}_{n}}\right)$ such that

$$
\begin{equation*}
\forall k \in\{1, \ldots, n\}, \quad B_{\mathbf{Y}_{k}} \subseteq S_{k} B_{\mathbf{X}_{k}} \quad \text { and } \quad \frac{\operatorname{iq}\left(B_{\mathbf{Y}_{k}}\right)}{\sqrt{m_{k}}}\left(\frac{\operatorname{vol}_{m_{k}}\left(B_{\mathbf{X}_{k}}\right)}{\operatorname{vol}_{m_{k}}\left(B_{\mathbf{Y}_{k}}\right)}\right)^{\frac{1}{m_{k}}} \leqslant \alpha \tag{64}
\end{equation*}
$$

Then, there exist a normed space $\mathbf{Y}=\left(\mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}},\|\cdot\|_{\mathbf{X}}\right)$ and $S \in S L\left(\mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}}\right)$ such that

$$
\begin{equation*}
B_{\mathbf{Y}} \subseteq S B_{\mathbf{X}} \quad \text { and } \quad \frac{\operatorname{iq}\left(B_{\mathbf{Y}}\right)}{\sqrt{m_{1}+\ldots+m_{n}}}\left(\frac{\operatorname{vol}_{m_{1}+\ldots+m_{n}}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{m_{1}+\ldots+m_{n}}\left(B_{\mathbf{Y}}\right)}\right)^{\frac{1}{m_{1}+\ldots+m_{n}}} \lesssim \alpha \tag{65}
\end{equation*}
$$

Since (64) with $\alpha=O(1)$ is immediate when $n_{0}=1$, Lemma 36 establishes Conjecture 10 for when $K$ is the unit ball of an unconditional normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$. This holds, in particular, for $\mathbf{X}=\ell_{p}^{n}$, though we will prove in Section 6.1 that the stronger conclusion of Conjecture 9 holds in this case (recall Remark 31. Lemma 36 also shows that Conjecture 10 holds for, say, $\mathbf{X}=\ell_{p}^{n}\left(\ell_{q}^{m}\right)$; we expect that the reasoning of Section 6.1 could be adapted to yield Conjecture 9 for these spaces as well, but we did not attempt to carry this out. Other spaces that satisfy (64) with $\alpha$ slowly growing will be presented in Section 1.6.2, upon their substitution into Lemma 36. more examples for which Conjecture 10 holds up to lower-order factors are obtained (of course, we are conjecturing here that it holds for any space).

Remark 37. Say that a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is in Cheeger position if

$$
\forall S \in \mathrm{SL}_{n}(\mathbb{R}), \quad \frac{\operatorname{vol}_{n-1}\left(\partial \operatorname{Ch} B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(\operatorname{Ch} B_{\mathbf{X}}\right)} \leqslant \frac{\operatorname{vol}_{n-1}(\partial \operatorname{ChSB\mathbf {X}})}{\operatorname{vol}_{n}\left(\operatorname{ChSB} B_{\mathbf{X}}\right)}
$$

Observe that if $\mathbf{X}$ is in Cheeger position, then its Cheeger space $\mathrm{Ch} \mathbf{X}$ is in minimum surface area position, namely, $\operatorname{vol}_{n-1}\left(\partial \operatorname{Ch} B_{\mathbf{X}}\right) \leqslant \operatorname{vol}_{n-1}\left(\partial S \operatorname{Ch} B_{\mathbf{X}}\right)$ for every $S \in S L_{n}(\mathbb{R})$. Indeed, $S \operatorname{Ch} B_{\mathbf{X}} \subseteq S B_{\mathbf{X}}$, so by the definition of the Cheeger body of $S B_{\mathbf{X}}$ we have $\operatorname{vol}_{n-1}\left(\partial S \operatorname{Ch} B_{\mathbf{X}}\right) / \operatorname{vol}_{n}\left(\operatorname{Ch} B_{\mathbf{X}}\right) \geqslant \operatorname{vol}_{n-1}\left(\partial \operatorname{ChS}_{\mathbf{X}}\right) / \operatorname{vol} l_{n}\left(\operatorname{Ch} S B_{\mathbf{X}}\right)$. At the same time, $\operatorname{vol}_{n-1}\left(\partial \operatorname{ChSB} B_{\mathbf{X}}\right) / \operatorname{vol}_{n}\left(\operatorname{ChSB} B_{\mathbf{X}}\right) \geqslant \operatorname{vol}_{n-1}\left(\partial \operatorname{Ch}_{\mathbf{X}}\right) / \operatorname{vol}_{n}\left(\operatorname{Ch} B_{\mathbf{X}}\right)$ by the definition of the Cheeger position, so $\operatorname{vol}_{n-1}\left(\partial S \operatorname{Ch} B_{\mathbf{X}}\right) \geqslant \operatorname{vol}_{n-1}\left(\partial \operatorname{Ch} B_{\mathbf{X}}\right)$. This shows that in the proof of the implication (2) $\Longrightarrow$ (1) of Proposition 32, if we worked with $L=\operatorname{ChSK}$, then there would be no need to introduce the additional linear transformation $T \in \mathrm{SL}_{n}(\mathbb{R})$. It would be worthwhile to study the Cheeger position for its own sake even if it weren't for its connection to reverse isoperimetry. In particular, we do not know if the converse of the above deduction holds, namely whether it is true that if $\operatorname{Ch} \mathbf{X}$ is in minimum surface area position, then $\mathbf{X}$ is in Cheeger position. We also do not know if the Cheeger position is unique up to orthogonal transformation (as is the case for the minimum surface area position [GP99]); we did not investigate these matters since they are not needed for the present purposes, but we expect that the characterisations of the Cheeger body in AC09 would be relevant here. One could also define that a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is in Dirichlet position if $\lambda(\mathbf{X}) \leqslant \lambda(S \mathbf{X})$ for every $S \in \mathrm{SL}_{n}(\mathbb{R})$. It is unclear how the Cheeger position relates to the Dirichlet position and it would be also worthwhile to study the Dirichlet position for its own sake. By (61), working with either the Cheeger position or the Dirichlet position would be equally valuable for the reverse isoperimetric questions in which we are interested here.
1.6.2. Symmetries and positions. Thus far we considered an arbitrary scalar product on an $n$-dimensional normed space through which we identified its underlying vector space structure with $\mathbb{R}^{n}$. However, the Lipschitz extension modulus is insufficiently understood for "very nice" normed spaces (including even the Euclidean space $\ell_{2}^{n}$ ) that belong to a natural class of normed spaces that have a canonical identification with $\mathbb{R}^{n}$. It therefore makes sense to first focus on this class.

For a finite dimensional normed space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$, let Isom $(\mathbf{X})$ be the group of all of the isometric automorphism of $\mathbf{X}$, i.e., all the linear operators $U: \mathbf{X} \rightarrow \mathbf{X}$ that satisfy $\|U x\|_{\mathbf{X}}=\|x\|_{\mathbf{X}}$ for every $x \in \mathbf{X}$. We will denote the Haar probability measure on the compact group Isom $(\mathbf{X})$ by $h_{\mathbf{X}}$.

Definition 38. We say that a finite dimensional normed space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ is canonically positioned if any two $\operatorname{lsom}(\mathbf{X})$-invariant scalar products on $\mathbf{X}$ are proportional to each other. In other words, if $\langle\cdot, \cdot\rangle: \mathbf{X} \times \mathbf{X} \rightarrow \mathbb{R}$ and $\langle\cdot, \cdot\rangle^{\prime}: \mathbf{X} \times \mathbf{X} \rightarrow \mathbb{R}$ are scalar products on $\mathbf{X}$ such that $\langle U x, U y\rangle=\langle x, y\rangle$ and $\langle U x, U y\rangle^{\prime}=\langle x, y\rangle^{\prime}$ for every $x, y \in \mathbf{X}$ and every $U \in \operatorname{Isom}(\mathbf{X})$, then there necessarily exists $\lambda \in \mathbb{R}$ such that $\langle\cdot, \cdot\rangle^{\prime}=\lambda\langle\cdot, \cdot\rangle$.

Any finite dimensional normed space $\mathbf{X}$ has at least one scalar product $\langle\cdot, \cdot \cdot\rangle: \mathbf{X} \times \mathbf{X} \rightarrow \mathbb{R}$ that is invariant under $\operatorname{Isom}(\mathbf{X})$, as seen e.g. by averaging any given scalar product $\langle\cdot, \cdot\rangle_{0}$ on $\mathbf{X}$ with respect $h_{\mathbf{X}}$, i.e., defining

$$
\forall x, y \in \mathbf{X}, \quad\langle x, y\rangle \stackrel{\text { def }}{=} \int_{\operatorname{Isom}(\mathbf{X})}\langle S x, S y\rangle_{0} \mathrm{~d} h_{\mathbf{X}}(S)
$$

Definition 38 concerns those spaces $\mathbf{X}$ for which such an invariant scalar product is unique up to rescaling, so there is (essentially, i.e., up to rescaling) no arbitrariness when we identify $\mathbf{X}$ with $\mathbb{R}^{\operatorname{dim}(\mathbf{X})}$.

Example 39. The class of $n$-dimensional canonically positioned normed spaces $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ includes those with a basis $e_{1}, \ldots, e_{n}$ such that for any distinct $i, j \in\{1, \ldots, n\}$ there are a permutation $\pi \in S_{n}$ with $\pi(i)=j$ and a sign vector $\varepsilon=\left(\varepsilon_{1}, \ldots, \varepsilon_{n}\right) \in\{-1,1\}^{n}$ with $\varepsilon_{i}=-\varepsilon_{j}$ such that $T_{\pi}, S_{\varepsilon} \in \operatorname{lsom}(\mathbf{X})$, where we denote $T_{\pi} x=\sum_{i=1}^{n} a_{\pi(i)} e_{i}$ and $S_{\varepsilon} x=\sum_{i=1}^{n} \varepsilon_{i} a_{i} e_{i}$ for $x=\sum_{i=1}^{n} a_{i} e_{i} \in \mathbf{X}$ with $a_{1}, \ldots, a_{n} \in \mathbb{R}$. Indeed, let $\langle\cdot, \cdot\rangle$ be a scalar product on $\mathbf{X}$ that is Isom( $\mathbf{X}$ )-invariant. For every distinct $i, j \in\{1, \ldots, n\}$, if $\pi \in S_{n}$ and $\varepsilon \in\{-1,1\}^{n}$ are as above, then $\left\langle e_{i}, e_{i}\right\rangle=\left\langle e_{\pi(i)}, e_{\pi(i)}\right\rangle=\left\langle e_{j}, e_{j}\right\rangle$ while $\left\langle e_{i}, e_{j}\right\rangle=\left\langle\varepsilon_{i} e_{i}, \varepsilon_{j} e_{j}\right\rangle=-\left\langle e_{i}, e_{j}\right\rangle$, so $\left\langle e_{i}, e_{j}\right\rangle=0$.

Example 39 covers all of the spaces for which we think that it is most pressing (given the current state of knowledge) to understand their Lipschitz extension modulus, including normed spaces $\left(\mathbf{E},\|\cdot\|_{\mathbf{E}}\right)$ that have a 1 -symmetric basis, i.e., a basis $e_{1}, \ldots, e_{n} \in \mathbf{E}$ such that $\left\|\sum_{i=1}^{n} \varepsilon_{i} a_{\pi(i)} e_{i}\right\|_{\mathbf{E}}=\left\|\sum_{i=1}^{n} a_{i} e_{i}\right\|_{\mathbf{E}}$ for every $(\varepsilon, \pi) \in\{-1,1\}^{n} \times S_{n}$. In particular, $\ell_{p}^{n}$, and more generally Orlicz and Lorentz spaces (see e.g. LTT77), are canonically positioned. We will use below the common convention that a normed space $\left(\mathbb{R}^{n},\|\cdot\|\right)$ is said to be symmetric if it is 1 -symmetric with respect to the standard (coordinate) basis $e_{1}, \ldots, e_{n}$ of $\mathbb{R}^{n}$.

Example 39 also includes matrix norms $\mathbf{X}=\left(\mathrm{M}_{n}(\mathbb{R}),\|\cdot\|_{\mathbf{X}}\right)$ that remain unchanged if one transposes a pair of rows or columns, or changes the sign of an entire row or a column, such as $S_{p}^{n}$. More generally, if $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ is a symmetric normed space, then its unitary ideal $\mathrm{S}_{\mathbf{E}}=\left(\mathrm{M}_{n}(\mathbb{R}),\|\cdot\|_{\mathrm{S}_{\mathrm{E}}}\right)$ is canonically positioned (see e.g. Bha97]), where for $T \in \mathrm{M}_{n}(\mathbb{R})$ one denotes its singular values by $s_{1}(T) \geqslant \ldots \geqslant s_{n}(T)$ and defines $\|T\|_{\mathrm{S}_{\mathbf{E}}}=\left\|\left(s_{1}(T), \ldots, s_{n}(T)\right)\right\|_{\mathbf{E}}$. More examples of such matrix norms are projective and injective tensor products (see e.g. $\overline{\text { Rya02| }}$ ) of symmetric spaces, where if $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ and $\mathbf{Y}=\left(\mathbb{R}^{m},\|\cdot\|_{\mathbf{Y}}\right)$ are normed spaces, then their projective tensor product $\mathbf{X} \hat{\otimes} \mathbf{Y}$ is the norm on $M_{n \times m}(\mathbb{R})=\mathbb{R}^{n} \otimes \mathbb{R}^{m}$ whose unit ball is the convex hull of $\left\{x \otimes y:(x, y) \in B_{\mathbf{X}} \times B_{\mathbf{Y}}\right\}$, and their injective tensor product $\mathbf{X} \otimes \check{\otimes}$ is the dual of $\mathbf{X}^{*} \hat{\otimes} \mathbf{Y}^{*}$ (equivalently, $\mathbf{X} \ddot{\otimes} \mathbf{Y}$ is isometric to the operator norm from $\mathbf{X}^{*}$ to $\mathbf{Y}$; see e.g. [DFS08, Section 1.1]).

Henceforth, when we will say that a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is canonically positioned it will always be tacitly assumed that the standard scalar product $\langle\cdot, \cdot\rangle$ on $\mathbb{R}^{n}$ is $\operatorname{Isom}(\mathbf{X})$-invariant, i.e., Isom $(\mathbf{X})$ is a subgroup of the orthogonal group $\mathrm{O}_{n} \subseteq \mathrm{M}_{n}(\mathbb{R})$. This is equivalent to the requirement that for every symmetric positive definite matrix $T \in \mathrm{M}_{n}(\mathbb{R})$, if $T U=U T$ for every $U \in \operatorname{Isom}(\mathbf{X})$, then there is $\lambda \in(0, \infty)$ such that $T=\lambda \operatorname{ld}_{n}$. Indeed, any scalar product $\langle\cdot, \cdot\rangle^{\prime}: \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}$ is of the form $\langle x, y\rangle^{\prime}=\langle T x, y\rangle$ for some symmetric positive definite $T \in \mathrm{M}_{n}(\mathbb{R})$ and all $x, y \in \mathbb{R}^{n}$, and using the Isom $(\mathbf{X})$-invariance of $\langle\cdot, \cdot\rangle$ we see that $\langle\cdot, \cdot\rangle^{\prime}$ is $\operatorname{lsom}(\mathbf{X})$-invariant if and only if $T$ commutes with all of the elements of Isom(X).

Remark 40. A symmetry assumption that is common in the literature is enough symmetries. A normed space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ is said [GG71] to have enough symmetries if any linear transformation $T: \mathbf{X} \rightarrow \mathbf{X}$ must be a scalar multiple of the identity if $T$ commutes with every element of $\operatorname{Isom}(\mathbf{X})$. By the above discussion, if $\mathbf{X}$ has enough symmetries, then $\mathbf{X}$ is canonically positioned. The converse implication does not hold, i.e., there exist normed spaces that are canonically positioned but do not have enough symmetries. For example, let $\operatorname{Rot}_{\pi / 2} \in \mathrm{O}_{2}$ be the rotation by 90 degrees and let $G$ be the subgroup of $\mathrm{O}_{2}$ that is generated by $\operatorname{Rot}_{\pi / 2}$. Thus, $G$ is cyclic of order 4 . Let $\mathbf{X}=\left(\mathbb{R}^{2},\|\cdot\|_{\mathbf{X}}\right)$ be a normed space with Isom $(\mathbf{X})=G$; the fact that there is such a normed space follows from the general result [GL79, Theorem 3.1] of Gordon and Loewy on existence of norms with a specified group of isometries, though in this particular case it is simple to construct such an example (e.g. the unit ball of $\mathbf{X}$ can be taken to be a suitable non-regular octagon). Since $\operatorname{Isom}(\mathbf{X})$ is Abelian, the matrix $\operatorname{Rot}_{\pi / 2}$ commutes with all of the elements of $\operatorname{Isom}(\mathbf{X})$ yet it is not a multiple of the identity matrix, so $\mathbf{X}$ does not have enough symmetries. Nevertheless, $\mathbf{X}$ is canonically positioned. Indeed, suppose that $T \in \mathrm{M}_{2}(\mathbb{R})$ is a symmetric matrix that commutes with $\operatorname{Rot}_{\pi / 2}$. Then, $\operatorname{Rot}_{\pi / 2}$ preserves any eigenspace of $T$, which means that any such eigenspace must be $\{0\}$ or $\mathbb{R}^{2}$. But $T$ is diagonalizable over $\mathbb{R}$, so it follows that $T=\lambda \mid \mathrm{d}_{2}$ for some $\lambda \in \mathbb{R}$. If $n$ is even, then one obtains such an $n$-dimensional example by considering $\ell_{\infty}^{n / 2}(\mathbf{X})$. However, a representation-theoretic argument due to Emmanuel Breuillard (private communication; details omitted) shows that if $n$ is odd, then any $n$-dimensional normed space has enough symmetries if and only if it is canonically positioned.

The following lemma is important for us even though it is an immediate consequence of the (major) theorem of [AC09] that the Cheeger body of a given convex body in $\mathbb{R}^{n}$ is unique (recall Section 1.6.1].
Lemma 41. Let $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathrm{X}}\right)$ be a normed space such that $\operatorname{Isom}(\mathbf{X}) \leqslant \mathrm{O}_{n}$ is a subgroup of the orthogonal group. Then the isometry group of its Cheeger space ChX satisfies

$$
\operatorname{Isom}(\mathrm{ChX}) \supseteq \operatorname{Isom}(\mathbf{X}) .
$$

Consequently, if $\mathbf{X}$ is canonically positioned, then also $\mathrm{Ch} \mathbf{X}$ is canonically positioned.
Proof. For any $U \in \operatorname{Isom}(\mathbf{X})$ we have $\operatorname{vol}_{n-1}\left(\partial U \operatorname{Ch}_{\mathbf{X}}\right) / \operatorname{vol}_{n}\left(U \operatorname{Ch} B_{\mathbf{X}}\right)=\operatorname{vol}_{n-1}\left(\partial \operatorname{Ch}_{\mathbf{X}}\right) / \operatorname{vol}_{n}\left(\operatorname{Ch} B_{\mathbf{X}}\right)$, since $U \in \mathrm{O}_{n}$, and also $U \operatorname{Ch} B_{\mathbf{X}} \subseteq U B_{\mathbf{X}}=B_{\mathbf{X}}$. Hence (by definition), $U \operatorname{Ch} B_{\mathbf{X}}$ is a Cheeger body of $B_{\mathbf{X}}$, so by the uniqueness of the Cheeger body we have $U \operatorname{Ch} B_{\mathbf{X}}=\operatorname{Ch} B_{\mathbf{X}}$. Therefore, $U \in \operatorname{lsom}(\operatorname{Ch} \mathbf{X})$.

The following corollary is a quick consequence of Lemma 41
Corollary 42. Let $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ be a symmetric normed space. Then, its Cheeger space $\mathrm{Ch} \mathbf{E}$ is also symmetric and there exists a (unique) symmetric normed space $\chi \mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\chi \mathbf{E}}\right)$ such that the Cheeger space of the unitary ideal $\mathrm{S}_{\mathbf{E}}$ is the unitary ideal of $\chi \mathbf{E}$, i.e., $\mathrm{ChS}_{\mathbf{E}}=\mathrm{S}_{\chi \mathbf{E}}$.

Proof. The assertion that ChE is symmetric coincides with the requirement that Isom( ChE ) contains the group $\{-1,1\}^{n} \rtimes S_{n}=\left\{T_{\varepsilon} S_{\pi}:(\varepsilon, \pi) \in\{-1,1\}^{n} \times S_{n}\right\} \leqslant \mathrm{O}_{n}$, where we recall the notation of Example 39 . Since we are assuming that $\operatorname{Isom}(\mathbf{E}) \supseteq\{-1,1\}^{n} \rtimes S_{n}$, this follows from Lemma 41 Next, for every $U, V \in \mathrm{O}_{n}$ define $R_{U, V}: \mathrm{M}_{n}(\mathbb{R}) \rightarrow \mathrm{M}_{n}(\mathbb{R})$ by $\left(A \in M_{n}(\mathbb{R})\right) \mapsto U A V$. Since Isom $\left(\mathrm{S}_{\mathrm{E}}\right)$ contains $\left\{R_{U, V}: U, V \in \mathrm{O}_{n}\right\}$, by Lemma 41 so does $\operatorname{Isom}\left(\mathrm{Ch}_{\mathbf{E}}\right)$. A normed space $\left(\mathrm{M}_{n}(\mathbb{R}),\|\cdot\|\right)$ that is invariant under $R_{U, V}$ for all $U, V \in \mathrm{O}_{n}$ is the unitary ideal of a symmetric normed space $\mathbf{F}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{F}}\right)$; see e.g. Bha97, Theorem IV.2.1]. This $\mathbf{F}$ is unique (consider the values of $\|\cdot\|_{S_{F}}$ on diagonal matrices), so we can introduce the notation $\mathbf{F}=\chi \mathbf{E}$.

The same reasoning as in the proof of Corollary 42 shows that if $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ is an unconditional normed space, then so is ChE. Thus, the space $\mathbf{Y}$ in Lemma 36 when $\mathbf{X}_{1}=\ldots=\mathbf{X}_{n}=\mathbb{R}$ that satisfies (65) can be taken to unconditional, as seen by an inspection of the proof of Lemma36(specifically, the operator $S$ in (65) that arises in this case is diagonal, so $S \mathbf{E}$ is also unconditional and we can take $\mathbf{Y}=\mathbf{C h} S \mathbf{E}$ ).

Problem 43. We associated above to every symmetric normed space $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ two symmetric normed spaces $\operatorname{ChE}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathrm{ChE}}\right)$ and $\chi \mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\chi \mathbf{E}}\right)$. It would be valuable to understand these auxiliary norms
on $\mathbb{R}^{n}$, and in particular how they relate to each other. By the definition of the Cheeger body, its convexity and uniqueness, ChE is the unique minimizer of the functional

$$
\begin{equation*}
\mathbf{F} \mapsto \frac{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{F}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{F}}\right)}=\frac{\int_{\partial B_{\mathbf{F}}} 1 \mathrm{~d} x}{\int_{B_{\mathbf{F}}} 1 \mathrm{~d} x} \tag{66}
\end{equation*}
$$

over all symmetric normed spaces $\mathbf{F}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{F}}\right)$ with $B_{\mathbf{F}} \subseteq B_{\mathbf{E}}$; denote the set of all such $\mathbf{F}$ by $\mathcal{S} y m\left(\subseteq B_{\mathbf{E}}\right)$. In contrast to $66, \chi \mathbf{E}$ is the unique minimizer of the functional

$$
\begin{equation*}
\mathbf{F} \mapsto \frac{\int_{\partial B_{\mathbf{F}}} \Pi_{1 \leqslant i<j \leqslant n}\left|x_{i}^{2}-x_{j}^{2}\right| \mathrm{d} x}{\int_{B_{\mathbf{F}}} \Pi_{1 \leqslant i<j \leqslant n}\left|x_{i}^{2}-x_{j}^{2}\right| \mathrm{d} x} \tag{67}
\end{equation*}
$$

over the same domain $\mathcal{S y m}\left(\subseteq B_{\mathbf{E}}\right)$. To justify (67), observe first that by Corollary 42 we know that $\chi \mathbf{E}$ is the unique minimizer of the following functional over $\mathcal{S} y m\left(\subseteq B_{\mathbf{E}}\right)$ :

$$
\begin{equation*}
\mathbf{F} \mapsto \frac{\operatorname{vol}_{n^{2}-1}\left(\partial B_{S_{\mathrm{F}}}\right)}{\operatorname{vol}_{n^{2}}\left(B_{S_{\mathrm{F}}}\right)}=\lim _{\varepsilon \rightarrow 0^{+}} \frac{\int\left(B_{\mathrm{S}_{\mathrm{F}}}+\varepsilon B_{\mathrm{S}_{2}^{n}}\right) \backslash B_{\mathrm{S}_{\mathrm{F}}}}{} 1 \mathrm{~d} x . \tag{68}
\end{equation*}
$$

We claim that for every $\mathbf{F} \in \mathcal{S} y m\left(\subseteq B_{\mathbf{E}}\right)$ and $\varepsilon>0$,

$$
\begin{equation*}
\left(B_{\mathrm{S}_{\mathbf{F}}}+\varepsilon B_{\mathrm{S}_{2}^{n}}\right) \backslash B_{\mathrm{S}_{\mathbf{F}}}=\left\{A \in M_{n}(\mathbb{R}): s(A) \stackrel{\text { def }}{=}\left(s_{1}(A), \ldots, s_{n}(A)\right) \in\left(B_{\mathbf{F}}+\varepsilon B_{\ell_{2}^{n}}\right) \backslash B_{\mathbf{F}}\right\} \tag{69}
\end{equation*}
$$

where we denote the singular values of $A \in \mathrm{M}_{n}(\mathbb{R})$ by $s_{1}(A) \geqslant \ldots \geqslant s_{n}(A)$. Indeed, if $A$ belongs to the right hand side of (69), then $\|s(A)\|_{\mathbf{F}}>1$ and $s(A)=x+y$ for $x, y \in \mathbb{R}^{n}$ that satisfy $\|x\|_{\mathbf{F}} \leqslant 1$ and $\|y\|_{\ell_{2}^{n}} \leqslant \varepsilon$. Write $A=U D V$, where $D \in \mathrm{M}_{n}(\mathbb{R})$ is the diagonal matrix whose diagonal is the vector $s(A) \in \mathbb{R}^{n}$, and $U, V \in \mathrm{O}_{n}$. Let $D(x), D(y) \in \mathrm{M}_{n}(\mathbb{R})$ be the diagonal matrices whose diagonals equal $x, y$, respectively. By noting that $\|A\|_{\mathrm{S}_{\mathrm{F}}}=\|s(A)\|_{\mathbf{F}}>1$ and $A=U D_{x} V+U D_{y} V$, where $\|U D(x) V\|_{\mathrm{S}_{\mathbf{F}}} \leqslant 1$ and $\|U D(y) V\|_{\mathrm{S}_{2}^{n}} \leqslant \varepsilon$, we conclude that $A$ belongs to the left hand side of (69). The reverse inclusion is less straightforward. If $A$ belongs to the left hand side of (69), then $\|A\|_{\mathrm{S}_{\mathbf{F}}}>1$ and $A=B+C$, where $B, C \in \mathrm{M}_{n}(\mathbb{R})$ satisfy $\|B\|_{\mathrm{S}_{\mathbf{F}}}=\|s(B)\|_{\mathbf{F}} \leqslant 1$ and $\|C\|_{S_{2}^{n}} \leqslant \varepsilon$. By an inequality of Mirsky Mir60] we have $\|s(A)-s(B)\|_{\ell_{2}^{n}} \leqslant\|A-B\|_{S_{2}^{n}}=\|C\|_{S_{2}^{n}} \leqslant \varepsilon$. Hence $s(A)=s(B)+(s(A)-s(B)) \in\left(B_{\mathbf{F}}+\varepsilon B_{\ell_{2}^{n}}\right) \backslash B_{\mathbf{F}}$, i.e., $A$ belongs to the right hand side of (69). With (69) established, since membership of a matrix $A$ in either $B_{\mathbf{F}}$ or $\left(B_{\mathbf{F}}+\varepsilon B_{\ell_{2}^{n}}\right) \backslash B_{\mathbf{F}}$ depends only on $s(A)$, by the Weyl integration formula [Wey39] (see [AGZ10, Proposition 4.1.3] for the formulation that we are using),

$$
\frac{\int_{\left(B_{\mathrm{S}_{\mathbf{F}}}+\varepsilon B_{\mathrm{S}_{2}^{n}}\right) \backslash B_{\mathrm{S}_{\mathbf{F}}}} 1 \mathrm{~d} x}{\int_{B_{\mathbf{F}}} 1 \mathrm{~d} x}=\frac{\int_{\left(B_{\mathbf{F}}+\varepsilon B_{\ell_{2}^{n}}\right) \backslash B_{\mathbf{F}}} \Pi_{1 \leqslant i<j \leqslant n}\left|x_{i}^{2}-x_{j}^{2}\right| \mathrm{d} x}{\int_{B_{\mathbf{F}}} \Pi_{1 \leqslant i<j \leqslant n}\left|x_{i}^{2}-x_{j}^{2}\right| \mathrm{d} x} .
$$

Thus (67) follows from (68). Analysing the functional in 67) seems nontrivial but likely tractable using ideas from random matrix theory. It would be especially interesting to treat the case $\mathbf{E}=\ell_{\infty}^{n}$. While we have a reasonably good understanding of the (isomorphic) geometry space $\mathrm{Ch} \ell_{\infty}^{n}$, its noncommutative counterpart $\chi \ell_{\infty}^{n}$ is still mysterious and understanding its geometry is closely related to Conjecture 10 (and likely also Conjecture 9 ) in the important special case of the operator norm $S_{\infty}^{n}$; see also Remark 171 .

If $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is canonically positioned and $\mu$ is a Borel measure on $\mathbb{R}^{n}$ that is Isom $(\mathbf{X})$-invariant, i.e., $\mu(U A)=\mu(A)$ for every $U \in \operatorname{Isom}(\mathbf{X})$ and every Borel subset $A \subseteq \mathbb{R}^{n}$, then consider the scalar product

$$
\forall x, y \in \mathbb{R}^{n}, \quad\langle x, y\rangle^{\prime} \stackrel{\text { def }}{=} \int_{\mathbb{R}^{n}}\langle x, z\rangle\langle y, z\rangle \mathrm{d} \mu(z) .
$$

For every $U \in \operatorname{lsom}(\mathbf{X})$ and $x, y \in \mathbb{R}^{n}$ we have

$$
\langle U x, U y\rangle^{\prime}=\int_{\mathbb{R}^{n}}\langle U x, z\rangle\langle U y, z\rangle \mathrm{d} \mu(z)=\int_{\mathbb{R}^{n}}\left\langle x, U^{-1} z\right\rangle\left\langle y, U^{-1} z\right\rangle \mathrm{d} \mu(z)=\int_{\mathbb{R}^{n}}\langle x, z\rangle\langle y, z\rangle \mathrm{d} \mu(z)=\langle x, y\rangle^{\prime},
$$

where the second step uses the $\operatorname{lsom}(\mathbf{X})$-invariance of $\langle\cdot, \cdot\rangle$, and the third step uses the Isom $(\mathbf{X})$-invariance of $\mu$. Hence $\langle x, y\rangle^{\prime}=\lambda\langle x, y\rangle$ for some $\lambda \in \mathbb{R}$ and every $x, y \in \mathbb{R}^{n}$. By considering the case $x=y$ of this identity and integrating over $x \in S^{n-1}$ one sees that necessarily $n \lambda=\int_{\mathbb{R}^{n}}\|z\|_{\ell_{2}^{n}}^{2} \mathrm{~d} \mu(z)$. Hence,

$$
\begin{equation*}
\forall x, y \in \mathbb{R}^{n}, \quad \int_{\mathbb{R}^{n}}\langle x, z\rangle\langle y, z\rangle \mathrm{d} \mu(z)=\frac{\int_{\mathbb{R}^{n}}\|z\|_{\ell_{2}^{n}}^{2} \mathrm{~d} \mu(z)}{n}\langle x, y\rangle . \tag{70}
\end{equation*}
$$

By establishing (70) we have shown that if $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a canonically positioned normed space, then any Isom(X)-invariant Borel measure on $\mathbb{R}^{n}$ is isotropic [GM00, BGVV14] (the converse also holds, i.e., $\mathbf{X}$ is canonically positioned if and only if every $\operatorname{lsom}(\mathbf{X})$-invariant Borel measure on $\mathbb{R}^{n}$ is isotropic). In particular, let $\sigma_{\mathbf{X}}$ be the measure on $S^{n-1}$ that is given by $\sigma_{\mathbf{X}}(A)=\operatorname{vol}_{n-1}\left(\left\{x \in \partial B_{\mathbf{X}}: N_{\mathbf{X}}(x) \in A\right\}\right)$ for every measurable $A \subseteq S^{n-1}$, where for $x \in \partial B_{\mathbf{X}}$ the vector $N_{\mathbf{X}}(x) \in S^{n-1}$ is the (almost-everywhere uniquely defined) unit outer normal to $\partial B_{\mathbf{X}}$ at $x$, i.e., recalling (30), we use the simpler notation $N_{B_{\mathbf{X}}}=N_{\mathbf{X}}$. In other words, $\sigma_{\mathbf{X}}$ is the image under the Gauss map of the $(n-1)$-dimensional Hausdorff measure on $\partial B_{\mathbf{X}}$. Then, $\sigma_{\mathbf{X}}$ is $\operatorname{Isom}(\mathbf{X})$-invariant because every $U \in \operatorname{Isom}(\mathbf{X})$ is an orthogonal transformation and $N_{\mathbf{X}} \circ U=U \circ N_{\mathbf{X}}$ almost everywhere on $\partial B_{\mathbf{X}}$. By [Pet61], this implies that $\mathbf{X}$ is in its minimum surface area position (recall the proof of $\operatorname{Proposition~32,~}$, so $\operatorname{MaxProj}\left(B_{\mathbf{X}}\right)=\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right) / \sqrt{n}$ by [GP99, Proposition 3.1].

The following corollary follows by substituting the above conclusion into Theorem 21 ,
Corollary 44. Suppose that $n \in \mathbb{N}$ and that $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ and $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ are two $n$-dimensional normed spaces. Suppose also that $\mathbf{Y}$ is canonically positioned and $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$. Then,

$$
\mathrm{e}(\mathbf{X}) \lesssim \frac{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{Y}}\right) \operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right) \sqrt{n}}
$$

The assumption in Corollary 44 that $\mathbf{Y}$ is canonically positioned can be replaced by the requirement $\operatorname{MaxProj}\left(B_{\mathbf{Y}}\right) \lesssim \operatorname{vol}_{n-1}\left(\partial B_{\mathbf{Y}}\right) / \sqrt{n}$, which is much less stringent. In particular, by GP99, Proposition 3.1] it is enough to assume here that $B_{\mathbf{Y}}$ is in its minimum surface area position; see also Section 6.2,

We will denote the John and Löwner ellipsoids of a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ by $\mathscr{J}_{\mathbf{x}}$ and $\mathcal{L}_{\mathbf{X}}$, respectively; see Hen12. Thus, $\mathcal{J}_{\mathbf{x}} \subseteq \mathbb{R}^{n}$ is the ellipsoid of maximum volume that is contained in $B_{\mathbf{X}}$ and $\mathcal{L}_{\mathbf{X}} \subseteq \mathbb{R}^{n}$ is the ellipsoid of minimum volume that contains $B_{\mathbf{X}}$. Both of these ellipsoids are unique (Joh48]. The volume ratio $\operatorname{vr}(\mathbf{X})$ of $\mathbf{X}$ and external volume ratio $\operatorname{evr}(\mathbf{X})$ of $\mathbf{X}$ are defined by

$$
\begin{equation*}
\operatorname{vr}(\mathbf{X}) \stackrel{\operatorname{def}}{=}\left(\frac{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(\mathcal{O}_{\mathbf{X}}\right)}\right)^{\frac{1}{n}} \quad \text { and } \quad \operatorname{evr}(\mathbf{X}) \stackrel{\operatorname{def}}{=}\left(\frac{\operatorname{vol}_{n}\left(\mathcal{L}_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}\right)^{\frac{1}{n}} \tag{71}
\end{equation*}
$$

By the Blaschke-Santaló inequality [Bla17, San49] and the Bourgain-Milman inequality [BM87],

$$
\begin{equation*}
\operatorname{evr}(\mathbf{X})=\operatorname{vr}\left(\mathbf{X}^{*}\right) \tag{72}
\end{equation*}
$$

By the above discussion, we can quickly deduce the following theorem that relates the Lipschitz extension modulus of a canonically positioned space to volumetric and spectral properties of its unit ball.
Theorem 45. Suppose that $n \in \mathbb{N}$ and that $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a canonically positioned normed space. Then,

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim \frac{\operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right)}{\sqrt{n}} \sqrt{\lambda(\mathbf{X})}=\operatorname{evr}(\mathbf{X}) \sqrt{\lambda(\mathbf{X}) \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{\frac{2}{n}}}=\operatorname{vr}\left(\mathbf{X}^{*}\right) \sqrt{\lambda(\mathbf{X}) \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{\frac{2}{n}}} \tag{73}
\end{equation*}
$$

In fact, the minimum of the right hand side of (53) over all the normed spaces $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ with $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$ is bounded above and below by universal constant multiples of $\operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right) \sqrt{\lambda(\mathbf{X}) / n}$.
Proof. By Lemma 41 the Cheeger space $\mathrm{Ch} \mathbf{X}$ is canonically positioned. So, by Corollary 44 with $\mathbf{Y}=\mathrm{ChX}$,

$$
\mathrm{e}(\mathbf{X}) \lesssim \frac{\operatorname{vol}_{n-1}\left(\partial \operatorname{Ch} B_{\mathbf{Y}}\right) \operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(\operatorname{Ch} B_{\mathbf{Y}}\right) \sqrt{n}} \stackrel{61}{\lesssim} \frac{\operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right)}{\sqrt{n}} \sqrt{\lambda(\mathbf{X})}
$$

This proves the first inequality in 73. The final equivalence in (73) is (72). To prove the rest of 73), let $r_{\text {min }}=\min \left\{r>0: r B_{\ell_{2}^{n}} \supseteq B_{\mathbf{X}}\right\}$ denote the radius of the circumscribing Euclidean ball of $B_{\mathbf{X}}$. We claim
that $r_{\min } B_{\ell_{2}^{n}}=\mathcal{L}_{\mathbf{X}}$. Indeed, for every $U \in \operatorname{Isom}(\mathbf{X}) \subseteq \mathrm{O}_{n}$ the ellipsoid $U \mathcal{L}_{\mathbf{X}}$ contains $B_{\mathbf{X}}$ and has the same volume as $\mathcal{L}_{\mathbf{X}}$, so because the minimum volume ellipsoid that contains $B_{\mathbf{X}}$ is unique [Joh48], it follows that $U \mathcal{L}_{\mathbf{X}}=\mathcal{L}_{\mathbf{X}}$. Hence, the scalar product that corresponds to $\mathcal{L}_{\mathbf{X}}$ is Isom $(\mathbf{X})$-invariant and since $\mathbf{X}$ is canonically positioned, this means that $\mathcal{L}_{\mathbf{X}}$ is a multiple of $B_{\ell_{2}^{n}}$. Now,

$$
\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{\frac{1}{n}} \operatorname{evr}(\mathbf{X}) \stackrel{(71)}{=} \operatorname{vol}_{n}\left(r_{\min } B_{\ell_{2}^{n}}\right)^{\frac{1}{n}}=\frac{r_{\text {min }}}{\sqrt{n}}=\frac{\operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right)}{2 \sqrt{n}} .
$$

The above reasoning shows that the minimum of the right hand side of (53) over all the normed spaces $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ with $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$ is at most a universal constant multiple of diam $\ell_{2}^{n}\left(B_{\mathbf{X}}\right) \sqrt{\lambda(\mathbf{X}) / n}$ (take $\left.\mathbf{Y}=\operatorname{Ch} \mathbf{X}\right)$. In the reverse direction, for any such $\mathbf{Y}$ by (54) with $L=B_{\mathbf{Y}}$ we have

$$
\frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \gtrsim \frac{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right) \sqrt{n}} \geqslant \frac{\operatorname{vol}_{n-1}\left(\partial \operatorname{Ch} B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(\operatorname{Ch} B_{\mathbf{X}}\right) \sqrt{n}} \stackrel{\sqrt[611]{2}}{\geqslant} \frac{2 \sqrt{\lambda(\mathbf{X})}}{\pi \sqrt{n}}
$$

where the penultimate step follows from the definition of the Cheeger body $\operatorname{Ch} B_{\mathbf{X}}$.
It is natural to expect that if $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a canonically positioned normed space, then in Conjecture 9 for $K=B_{\mathbf{X}}$ holds with $S$ the identity matrix and with $L$ being the unit ball of a canonically positioned normed space. We formulate this refined special case of Conjecture 9 as the following conjecture.

Conjecture 46. Fix $n \in \mathbb{N}$ and a canonically positioned normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$. Then, there exists a canonically positioned normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ with $\|\cdot\|_{\mathbf{Y}}=\|\cdot\|_{\mathbf{X}}$ and $\mathrm{iq}\left(B_{\mathbf{Y}}\right) \lesssim \sqrt{n}$.

Theorem 47 below shows that Conjecture 46 holds if $\mathbf{X}=\ell_{p}^{n}$ for any $p \geqslant 1$ and infinitely many dimensions $n \in \mathbb{N}$; specifically, it holds if $n$ satisfies the mild arithmetic (divisibility) requirement (74) below. An obvious question that this leaves is to prove Conjecture 46 for $\mathbf{X}=\ell_{p}^{n}$ and arbitrary $(p, n) \in[1, \infty] \times \mathbb{N}$. We expect that this question is tractable by (likely nontrivially) adapting the approach herein, but we did not make a major effort to do so since obtaining Conjecture 46 for such a dense set of dimensions $n$ suffices for our purposes (the bi-Lipschitz invariants that we consider can be estimated from above for any $n \in \mathbb{N}$ since the requirement (74) holds for some $N \in \mathbb{N} \cap[n, O(n)]$ and $\ell_{p}^{n}$ embeds isometrically into $\ell_{p}^{N}$ ). In Section 6, we will prove Theorem 47, and deduce Theorem 24 from it. Recall Remark 31 , which explains that Conjecture 9 when $K$ is the unit ball of $\ell_{p}^{n}$ follows (with $S$ the identity matrix) from Theorem 24 . Thus, we do know that a body $L$ as in Conjecture 9 exists for all the possible choices of $p \geqslant 1$ and $n \in \mathbb{N}$, and $(74)$ is only relevant to ensure that $L$ is the unit ball of a canonically positioned normed space.

Theorem 47. Fix $n \in \mathbb{N}$ and $p \geqslant 1$. Conjecture 46 holds for $\mathbf{X}=\ell_{p}^{n}$ if the following condition is satisfied.

$$
\begin{equation*}
\exists m \in \mathbb{N}, \quad m \mid n \quad \text { and } \quad \max \{p, 2\} \leqslant m \leqslant e^{p} \tag{74}
\end{equation*}
$$

The following conjecture is a variant of Conjecture 11 .
Conjecture 48. Fix $n \in \mathbb{N}$ and a canonically positioned normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$. There exists a normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ with $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$ yet $\sqrt[n]{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \gtrsim \sqrt[n]{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}$ such that $\mathrm{iq}\left(B_{\mathbf{Y}}\right) \lesssim \sqrt{n}$.

Conjecture 46 requires $\mathbf{Y}$ to be canonically positioned while Conjecture 48 does not. The reason for this is that if any normed space $\mathbf{Y}$ satisfies the conclusion of Conjecture 48 , then also the Cheeger space ChX of $\mathbf{X}$ satisfies it (this is so because the convex body $L$ that minimizes the second quantity in (60) is, by definition, the Cheeger body of $K=B_{\mathbf{X}}$ ), and by Lemma 41 the Cheeger space of $\mathbf{X}$ inherits from $\mathbf{X}$ the property of being canonically positioned. This use of the uniqueness of the Cheeger body will be important below. By 62, Conjecture 48 is equivalent to the following symmetric version of Conjecture 35 .
Conjecture 49. If $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a canonically positioned normed space, then $\lambda(\mathbf{X}) \operatorname{vol}\left(B_{\mathbf{X}}\right)^{\frac{2}{n}}=n$.
The following corollary is a substitution of Conjecture 49 into Theorem 45 .

Corollary 50. If Conjecture 48 (equivalently, Conjecture 49) holds for a canonically positioned normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$, then the right hand side of (53) when $\mathbf{Y}=\operatorname{ChX}$ is $O(\operatorname{evr}(\mathbf{X}) \sqrt{n})$. Consequently,

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim \mathrm{evr}(\mathbf{X}) \sqrt{n}=\operatorname{vr}\left(\mathbf{X}^{*}\right) \sqrt{n} . \tag{75}
\end{equation*}
$$

It is worthwhile to note that by [Bal89], the rightmost quantity in (75) is maximized (over all possible $n$-dimensional normed spaces) when $\mathbf{X}=\ell_{1}^{n}$, in which case we have $\operatorname{evr}\left(\ell_{1}^{n}\right) \sqrt{n}=n$.
Remark 51. We currently do not have any example of a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ for which (75) provably does not hold. If (75) were true in general, or even if it were true for a restricted class of normed spaces that is affine invariant and closed under direct sums, such as spaces that embed into $\ell_{1}$ with distortion $O(1)$, then it would be an excellent result. When one leaves the realm of canonically positioned spaces, (75) acquires a self-improving property $]^{8}$ as follows. Suppose that $\mathbf{X}$ is in Löwner position, i.e., $\mathcal{L}_{\mathbf{X}}=B_{\ell_{2}^{n}}$. Fix $m \in \mathbb{N}$ and consider the $(n+m)$-dimensional space $\mathbf{X}^{\prime}=\mathbf{X} \oplus_{\infty} \ell_{2}^{m}$. If (75) holds for $\mathbf{X}^{\prime}$, then

$$
\begin{align*}
\mathrm{e}(\mathbf{X}) \leqslant \mathrm{e}\left(\mathbf{X}^{\prime}\right) & \lesssim \operatorname{evr}\left(\mathbf{X}^{\prime}\right) \sqrt{\operatorname{dim}\left(\mathbf{X}^{\prime}\right)} \lesssim\left(\frac{\operatorname{vol}_{n+m}\left(B_{\ell_{2}^{n+m}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) \operatorname{vol}_{m}\left(B_{\left.\ell_{2}^{m}\right)}\right.}\right)^{\frac{1}{n+m}} \sqrt{n+m}  \tag{76}\\
& =\left(\frac{\operatorname{vol}_{n}\left(\mathcal{L}_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}\right)^{\frac{1}{n+m}}\left(\frac{\operatorname{vol}_{n+m}\left(B_{\ell_{2}^{n+m}}\right)}{\operatorname{vol}_{n}\left(\ell_{2}^{n}\right) \operatorname{vol}_{m}\left(B_{\ell_{2}^{m}}\right)}\right)^{\frac{1}{n+m}} \sqrt{n+m}=\operatorname{evr}\left(\mathbf{X} \frac{n}{n^{n+m}} n^{\frac{n}{2(n+m)}} m^{\frac{m}{2(n+m)}} .\right.
\end{align*}
$$

The value of $m$ that minimizes the right hand side of (76) is $m=n \log (\operatorname{evr}(\mathbf{X})+1)$, for which (76) becomes

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim \sqrt{n \log (\operatorname{evr}(\mathbf{X})+1)} \tag{77}
\end{equation*}
$$

As $\operatorname{evr}(\mathbf{X}) \leqslant \sqrt{n}$ by John's theorem, (77) gives $\mathrm{e}(\mathbf{X}) \lesssim \sqrt{n \log n}$, which would be an improvement of [JLS86]. Also, by (9) the bound (77) gives $\mathrm{e}(\mathbf{X}) \lesssim \sqrt{n \log \left(C_{2}(\mathbf{X})+1\right)}$, which is better than the conjectural bound (10). Here and throughout, for $1 \leqslant p \leqslant 2 \leqslant q$ the (Gaussian) type- $p$ and cotype- $q$ constants [MP76] of a Banach space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$, denoted $T_{p}(\mathbf{X})$ and $C_{q}(\mathbf{X})$, respectively, are the infimum over those $T \in[1, \infty]$ and $C \in[1, \infty]$, respectively, for which the following inequalities hold for every $m \in \mathbb{N}$ and every $x_{1}, \ldots, x_{m} \in \mathbf{X}$, where the expectation is with respect to i.i.d. standard Gaussian random variables $\mathrm{g}_{1}, \ldots, \mathrm{~g}_{m}$.

$$
\begin{equation*}
\frac{1}{C}\left(\sum_{j=1}^{m}\left\|x_{j}\right\|_{\mathbf{X}}^{q}\right)^{\frac{1}{q}} \leqslant\left(\mathbb{E}\left[\left\|\sum_{j=1}^{m} \mathrm{~g}_{j} x_{j}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \leqslant T\left(\sum_{j=1}^{m}\left\|x_{j}\right\|_{\mathbf{X}}^{p}\right)^{\frac{1}{p}} \tag{78}
\end{equation*}
$$

This observation indicates that it might be too optimistic to expect that (75) holds in full generality, but it would be very interesting to understand the extent to which it does. Obvious potential counterexamples are $\ell_{1}^{n} \oplus \ell_{2}^{m}$; if (75) holds for these spaces, then $\left(\ell_{1}^{n}\right) \lesssim \sqrt{n \log n}$ by the above reasoning (with $m=n \log n$ ), which would be a big achievement because the best-known bound remains e $\left(\ell_{1}^{n}\right) \lesssim n$ from (JLS86].

Lemma 52 below, whose proof appears in Section 6.1, shows that Conjecture 48 holds for a class of normed space that includes any normed spaces with a 1-symmetric basis, as well as, say, $\ell_{p}^{n}\left(\ell_{q}^{m}\right)$ for any $n, m \in \mathbb{N}$ and $p, q \geqslant 1$. Other (related) examples of such spaces arise from Lemma 150 below.
Lemma 52. Let $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ be an unconditional normed space. Suppose that for any $j, k \in\{1, \ldots, n\}$ there is a permutation $\pi \in S_{n}$ with $\pi(j)=k$ such that $\left\|\sum_{i=1}^{n} a_{\pi(i)} e_{i}\right\|_{\mathbf{X}}=\left\|\sum_{i=1}^{n} a_{i} e_{i}\right\|_{\mathbf{X}}$ for every $a_{1}, \ldots, a_{n} \in \mathbb{R}$. Then, Conjecture 48 holds for $\mathbf{X}$. Therefore, we have $\lambda(\mathbf{X}) \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{2 / n}=n$ and $\mathrm{e}(\mathbf{X}) \lesssim \operatorname{evr}(\mathbf{X}) \sqrt{n}$.

By [STJ80, Theorem 2.1], any unconditional normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{x}}\right)$ satisfies $\operatorname{vr}(\mathbf{X}) \lesssim C_{2}(\mathbf{X}) \sqrt{n}$, where $C_{2}(\mathbf{X})$ is the cotype-2 constant of $\mathbf{X}$ (this is an earlier special case of (9) in which the logarithmic term is known to be redundant). Hence, if $\mathbf{X}$ satisfies the assumptions of Lemma52, then we know that

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim C_{2}\left(\mathbf{X}^{*}\right) \sqrt{n} . \tag{79}
\end{equation*}
$$

[^4]By combining [Bal91c Theorem 6] and (72), for any $p \in[1, \infty]$, if a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is isometric to a quotient of $L_{p}$ (equivalently, the dual of $\mathbf{X}$ is isometric to a subspace of $L_{p /(p-1)}$ ), then

$$
\operatorname{evr}(\mathbf{X}) \lesssim \operatorname{evr}\left(\ell_{\frac{p}{p-1}}^{n}\right)=\min \left\{n^{\frac{1}{p}-\frac{1}{2}}, 1\right\}
$$

Consequently, if $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ satisfies the assumptions of Lemma 52 and is also a quotient of $L_{p}$, then

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim n^{\max \left\{\frac{1}{2}, \frac{1}{p}\right\}} \tag{80}
\end{equation*}
$$

Both (79) and (80) are generalizations of Theorem 18 .
Lemma 53 below, whose proof appears in Section 6.3 shows that the unitary ideal of any $n$-dimensional normed space with a 1-symmetric basis (in particular, any Schatten-von Neumann trace class), satisfies Conjecture 48 up to a factor of $O(\sqrt{\log n})$. Upon its substitution into Lemma 150 below, more such examples are obtained.

Lemma 53. Let $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ be a symmetric normed space. Conjecture 48 holds up to lower order factors for its unitary ideal $\mathrm{S}_{\mathrm{E}}$. More precisely, there is a normed space $\mathbf{Y}=\left(\mathrm{M}_{n}(\mathbb{R}),\|\cdot\| \mathbf{Y}\right)$ such that $B_{\mathbf{Y}} \subseteq B_{\mathrm{S}_{\mathrm{E}}}$ and

$$
\begin{equation*}
\operatorname{vol}_{n^{2}}\left(B_{\mathbf{Y}}\right)^{\frac{1}{n^{2}}}=\operatorname{vol}_{n^{2}}\left(B_{\mathrm{S}_{\mathrm{E}}}\right)^{\frac{1}{n^{2}}} \quad \text { and } \quad n \lesssim \operatorname{iq}\left(B_{\mathbf{Y}}\right) \lesssim n \sqrt{\log n} \tag{81}
\end{equation*}
$$

Therefore, we have

$$
n^{2} \lesssim \lambda\left(\mathrm{~S}_{\mathbf{E}}\right) \operatorname{vol}_{n^{2}}\left(B_{\mathrm{S}_{\mathbf{E}}}\right)^{\frac{2}{n^{2}}} \lesssim n^{2} \log n \quad \text { and } \quad \mathrm{e}\left(\mathrm{~S}_{\mathbf{E}}\right) \lesssim \operatorname{evr}\left(\mathrm{S}_{\mathbf{E}}\right) n=\operatorname{evr}(\mathbf{E}) n
$$

For the final assertion of Lemma 53 , the fact that $\operatorname{evr}\left(\mathrm{S}_{\mathbf{E}}\right)=\operatorname{evr}(\mathbf{E})$ follows by combining Proposition 2.2 in [Sch82], which states that $\operatorname{vr}\left(\mathrm{S}_{\mathbf{E}}\right)=\operatorname{vr}(\mathbf{E})$, with (72) and the duality $\mathrm{S}_{\mathbf{E}}^{*}=\mathrm{S}_{\mathbf{E}^{*}}$ (e.g. Sim79, Theorem 1.17]).

The proof of Lemma 53 also shows (see Remark 171 below) that if we could prove Conjecture 48 for $S_{\infty}^{n}$, then it would follow that $\mathrm{S}_{\mathbf{E}}$ satisfies Conjecture 48 for any symmetric normed space $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$, i.e., the logarithmic factor in (81) could be replaced by a universal constant.

By substituting Lemma 53 into Corollary 50 and using volume ratio computations of Schütt [Sch82], we will derive in Section 6.3 the following proposition.

Proposition 54. If $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ is a symmetric normed space, then

$$
\mathrm{e}(\mathbf{E}) \lesssim \operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{E}}\right)\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}} \quad \text { and } \quad \mathrm{e}\left(\mathrm{~S}_{\mathbf{E}}\right) \lesssim \operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{E}}\right)\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}} \sqrt{n \log n}
$$

The following remark sketches an alternative approach towards Conjecture 9 when $K$ is the hypercube $[-1,1]^{n}$ that differs from how we will prove Theorem 24 . It yields the desired result up to a lower order factor that grows extremely slowly; specifically, it constructs an origin-symmetric convex body $L \subseteq[-1,1]^{n}$ for which $[-1,1]^{n} \subseteq \exp \left(O\left(\log ^{*} n\right)\right) L$ and $\operatorname{iq}(L)=\exp \left(O\left(\log ^{*} n\right)\right.$. Here, for each $x \geqslant 1$ the quantity $\log ^{*} x$ is defined to be the $k \in \mathbb{N}$ such that tower $(k-1) \leqslant x<\operatorname{tower}(k)$ for the sequence $\{\operatorname{tower}(i)\}_{i=0}^{\infty}$ that is defined by tower $(0)=1$ and $\operatorname{tower}(i+1)=\exp (\operatorname{tower}(i))$. We think that this approach is worthwhile to describe despite the fact that it falls slightly short of fully establishing Conjecture 9 for $[-1,1]^{n}$ due to its flexibility that could be used for other purposes, as well as due to its intrinsic interest.

Remark 55. Fix $n \in \mathbb{N}$ and $q \geqslant 1$. Since the $n^{\prime}$ th root of the volume of the unit ball of $\ell_{q}^{n}$ is of order $n^{-1 / q}$ and $\ell_{q}^{n}$ is in minimum surface area position, we can restate (42) as

$$
\begin{equation*}
\operatorname{iq}\left(B_{\ell_{q}^{n}}\right)=\min \{\sqrt{q n}, n\} \tag{82}
\end{equation*}
$$

In particular, for $\mathbf{Y}=\ell_{q}^{n}$ with $q=\log n$, we have $\|\cdot\|_{\mathbf{Y}}=\|\cdot\|_{\ell_{\infty}^{n}}$ and iq $(\mathbf{Y}) \lesssim \sqrt{n \log n}$, which already comes close to the conclusion of Conjecture 9 . We can do better using the following evaluation of the isoperimetric quotient of the unit ball of $\ell_{p}^{n}\left(\ell_{q}^{m}\right)$, which holds for every $n, m \in \mathbb{N}$ and $p, q \geqslant 1$.

$$
\operatorname{iq}\left(B_{\ell_{p}^{n}\left(\ell_{q}^{m}\right)}\right)= \begin{cases}n m & m \leqslant \min \left\{\frac{p}{n}, q\right\},  \tag{83}\\ n \sqrt{q m} & q \leqslant m \leqslant \frac{p}{n}, \\ \sqrt{p n m} & \frac{p}{n} \leqslant m \leqslant \min \{p, q\}, \\ \sqrt{p q n} & \max \left\{\frac{p}{n}, q\right\} \leqslant m \leqslant p, \\ m \sqrt{n} & p \leqslant m \leqslant q, \\ \sqrt{q n m} & m \geqslant \max \{p, q\} .\end{cases}
$$

We will prove (83) in Section 6. Note that when $m=1$ this yields (82). The case $n=m$ of (83) is equivalent to (49) since $\ell_{p}^{n}\left(\ell_{q}^{m}\right)$ is canonically positioned (it belongs to the class of spaces in Example 39) and using a simple evaluation of the volume of its unit ball (see (316) below). The range of (83) that is most pertinent for the present context is $m \geqslant \max \{p, q\}$, which has the feature that the factor that multiplies the quantity

$$
\sqrt{n m}=\sqrt{\operatorname{dim}\left(\ell_{p}^{n}\left(\ell_{q}^{m}\right)\right)}
$$

is $O(\sqrt{q})$ and there is no dependence on $p$. This can be used as follows. Suppose that $n=a b$ for $a, b \in \mathbb{N}$ satisfying $a=n / \log n$ and $b=\log n$. Identify $\ell_{\infty}^{n}$ with $\ell_{\infty}^{a}\left(\ell_{\infty}^{b}\right)$. If we set $\mathbf{Y}=\ell_{p}^{a}\left(\ell_{q}^{b}\right)$ for $p=\log a=\log n$ and $q=\log b=\log \log n$, then $\|\cdot\|_{\mathbf{Y}}=\|\cdot\|_{\ell_{\infty}^{n}}$, while iq $\left(B_{\mathbf{Y}}\right)=\sqrt{n \log \log n}$ by (83). By iterating we get that for infinitely many $n \in \mathbb{N}$ there is a normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ for which $\|\cdot\|_{\mathbf{Y}} \leqslant\|\cdot\|_{\ell_{\infty}^{n}} \leqslant \exp \left(O\left(\log ^{*} n\right)\right)\|\cdot\|_{\mathbf{Y}}$ and $\operatorname{iq}\left(B_{\mathbf{Y}}\right)=\exp \left(O\left(\log ^{*} n\right)\right)$. Even though the set of $n \in \mathbb{N}$ for which this works is not all of $\mathbb{N}$, it is quite dense in $\mathbb{N}$ per Lemma 162 below. This will allow us to deduce that a space $\mathbf{Y}$ with the above properties exists for every $n \in \mathbb{N}$; see Section 6.1 for the details.
Remark 56. Recalling Remark 37, Conjecture 10 is equivalent to the assertion that if a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is in Cheeger position, then $\operatorname{vol}_{n}\left(\operatorname{Ch} B_{\mathbf{X}}\right)^{1 / n} \gtrsim \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{1 / n}$ and $\mathrm{iq}\left(\operatorname{Ch} B_{\mathbf{X}}\right) \lesssim \sqrt{n}$. Since $\operatorname{ChX}$ is in minimum surface area position when $\mathbf{X}$ is in Cheeger position (as explained in Remark 37), the proof of Proposition 32 shows that Conjecture 10 implies that if $\mathbf{X}$ is in Cheeger position, then

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \lesssim \frac{\operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{\frac{1}{n}}} \tag{84}
\end{equation*}
$$

In fact, the right hand side of (53) is at most the right hand side of (84) for a suitable choice of normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$, specifically for $\mathbf{Y}=\mathbf{C h X}$. The discussion in Section 1.6 .2 was about establishing (84) when $\mathbf{X}$ is canonically positioned (conceivably that assumption implies that $\mathbf{X}$ is in Cheeger position or close to it, which would be a worthwhile to prove, if true). Even though, as we explained earlier, given the current state of knowledge, understanding the Lipschitz extension problem for canonically positioned spaces is the most pressing issue for future research, it would be very interesting to study if (84) holds in other situations. For examples, we pose the following two natural questions.

Question 57. Does (84) hold if the normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathrm{X}}\right)$ is in minimum surface area position?
The extent to which $\Pi \mathbf{X}$ is close to being in minimum surface area position when $\mathbf{X}$ is in minimum surface area position seems to be unknown. Therefore, the connection between Question 58 below and Question 57 is unclear, but even if there is no formal link between these two questions, both are natural next steps beyond the setting of canonically positioned normed spaces.

Question 58 . Let $\mathbf{Z}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Z}}\right)$ be a normed space in minimum surface area position. Does (84) hold for the normed space $\mathbf{X}=\Pi \mathbf{Z}$ whose unit ball is the projection body of $B_{\mathbf{X}}$ ?

If $\mathbf{Z}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Z}}\right)$ is a normed space in minimum surface area position, then

$$
\begin{equation*}
\frac{\operatorname{diam}_{\ell_{2}^{n}\left(\Pi B_{\mathbf{Z}}\right)}}{\operatorname{vol}_{n}\left(\Pi B_{\mathbf{Z}}\right)^{\frac{1}{n}}}=\sqrt{n} . \tag{85}
\end{equation*}
$$

Indeed, because $\mathbf{Z}$ is in minimum surface area position, we have $\operatorname{vol}_{n}\left(\Pi B_{\mathbf{Z}}\right)^{1 / n}=\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{Z}}\right) / n$ by [GP99, Corollary 3.4], and also $\operatorname{MaxProj}\left(B_{\mathbf{Z}}\right)=\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{Z}}\right) / \sqrt{n}$ by combining [GP99, Proposition 3.1] and 54]. We can therefore justify (85) using these results from [GP99] and duality as follows.

$$
\frac{\operatorname{diam}_{\ell_{2}^{n}}\left(\Pi B_{\mathbf{Z}}\right)}{\operatorname{vol}_{n}\left(\Pi B_{\mathbf{Z}}\right)^{\frac{1}{n}}}=\frac{n\left\|\operatorname{ld}_{n}\right\|_{\Pi \mathbf{Z} \rightarrow \ell_{2}^{n}}}{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{Z}}\right)}=\frac{n\left\|\operatorname{ld}_{n}\right\|_{\ell_{2}^{n} \rightarrow \Pi^{*} \mathbf{Z}}}{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{Z}}\right)}=\frac{n \max _{z \in S^{n-1}}\|z\|_{\Pi^{*} \mathbf{Z}}}{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{Z}}\right)} \stackrel{30}{=} \frac{n \operatorname{MaxProj}\left(B_{\mathbf{Z}}\right)}{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{Z}}\right)}=\sqrt{n}
$$

By this observation, a positive answer to Question 58 would show that $\mathrm{e}(\Pi \mathbf{Z}) \lesssim \sqrt{n}$ for any normed space $\mathbf{Z}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Z}}\right)$. Indeed, if we take $S \in \mathrm{SL}_{n}(\mathbb{R})$ such that $S \mathbf{Z}$ is in minimum surface area position, then by [Pet67] we know that $\Pi \mathbf{Z}$ and $\Pi S \mathbf{Z}$ are isometric, so $e(\Pi \mathbf{Z})=\mathrm{e}(\Pi S \mathbf{Z})$. As the class of projection bodies coincides with the class of zonoids [Bol69, SW83], which coincides with the class of convex bodies whose polar is the unit ball of a subspace of $L_{1}$, we have thus shown that a positive answer to Question 58 would imply the following conjecture (which would simultaneously improve (23) and generalize Theorem 18).

Conjecture 59 . For any normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ we have $\mathrm{e}(\mathbf{X}) \lesssim \mathrm{c}_{L_{1}}\left(\mathbf{X}^{*}\right) \sqrt{n}$.
Note that Conjecture 59 is consistent with the estimate $\mathrm{e}(\mathbf{X}) \lesssim \mathrm{evr}(\mathbf{X}) \sqrt{n}$ that has been arising thus far. Indeed, if $\mathbf{X}^{*}$ is isometric to a subspace of $L_{1}$ (it suffices to consider only this case in Conjecture 59 by a well-known differentiation argument; see e.g. [BL00, Corollary 7.10]), then we have the bound evr(X) $\lesssim 1$ which can be seen to hold by combining (72) with (9), since $C_{2}\left(\mathbf{X}^{*}\right) \leqslant C_{2}\left(L_{1}\right) \lesssim 11^{9}$

Relating $e(\mathbf{X})$ to $\operatorname{evr}(\mathbf{X})$ is valuable since the Lipschitz extension modulus is for the most part shrouded in mystery, while the literature contains extensive knowledge on volume ratios (we have already seen several examples of such consequences above, and we will derive more later). Section 6.3 contains examples of volume ratio evaluations for various canonically positioned normed spaces. Through their substitution into Corollary50, they illustrate how our work yields a range of new Lipschitz extension results, some of which are currently conjectural because they hold assuming Conjecture 48 for the respective spaces; specifically, consider the Lipschitz extension bounds that correspond to using (14) and (15) with [LN05].
1.6.3. Intersection with a Euclidean ball. Fix an integer $n \geqslant 2$ and a canonically positioned normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$. A natural first attempt to prove Conjecture 48 for $\mathbf{X}$ is to consider the normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ such that $B_{\mathbf{Y}}=B_{\mathbf{X}} \cap r B_{\ell_{2}^{n}}$ for a suitably chosen $r>0$ (equivalently, $\|x\|_{\mathbf{Y}}=\max \left\{\|x\|_{\mathbf{X}},\|x\|_{\ell_{2}^{n}} / r\right\}$ for every $x \in \mathbb{R}^{n}$ ). However, we checked with $G$. Schechtman that this fails even when $\mathbf{X}=\ell_{\infty}^{n}$. Specifically, if the $n$ 'th root of the volume of $B_{\ell_{\infty}^{n}} \cap\left(r B_{\ell_{2}^{n}}\right)$ is at least a universal constant, then necessarily $r \gtrsim \sqrt{n}$, but

$$
\begin{equation*}
\forall s>0, \quad \operatorname{iq}\left(B_{\ell_{\infty}^{n}} \cap\left(s \sqrt{n} B_{\ell_{2}^{n}}\right)\right) \gtrsim s n . \tag{86}
\end{equation*}
$$

A justification of (86) appears in Section 7 below. In terms of the quantification (60) of Conjecture 48 that is pertinent to the applications that we study herein, we will also show in Section 7 that

$$
\begin{equation*}
\min _{r>0} \frac{\operatorname{iq}\left(B_{\ell_{\infty}^{n}} \cap\left(r B_{\ell_{2}^{n}}\right)\right)}{\sqrt{n}}\left(\frac{\operatorname{vol}_{n}\left(B_{\ell_{\infty}^{n}}\right)}{\operatorname{vol}_{n}\left(B_{\ell_{\infty}^{n}} \cap\left(r B_{\ell_{2}^{n}}\right)\right)}\right)^{\frac{1}{n}}=\sqrt{\log n} \tag{87}
\end{equation*}
$$

where the minimum in the right hand side of (87) is attained at some $r>0$ that satisfies $r=\sqrt{n / \log n}$.
Even though the above bounds demonstrate that it is impossible to resolve Conjecture 48 by intersecting with a Euclidean ball, this approach cannot fail by more than a lower-order factor; the reasoning that proves this assertion was shown to us by B. Klartag and E. Milman in unpublished private communication that is explained with their permission in Section 7 . Specifically, we have the following proposition.

[^5]Proposition 60. For any normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ there exist a matrix $S \in \mathrm{SL}_{n}(\mathbb{R})$ and a radius $r>0$ such that for $L=\left(S B_{\mathbf{X}}\right) \cap\left(r B_{\ell_{2}^{n}}\right) \subseteq S B_{\mathbf{X}}$ we have $\mathrm{iq}(L) \lesssim \sqrt{n}$ and $\sqrt[n]{\operatorname{vol}_{n}(L)} \gtrsim \sqrt[n]{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} / K(\mathbf{X})$, where $K(\mathbf{X})$ is the $K$-convexity constant of $\mathbf{X}$. If $\mathbf{X}$ is canonically positioned, then this holds when $S$ is the identity matrix.

For Proposition 60, the $K$-convexity constant of $\mathbf{X}$ is an isomorphic invariant that was introduced by Maurey and Pisier [MP76] we defer recalling its definition to Section 7 since for the discussion here it suffices to state the following bounds that relate $K(\mathbf{X})$ to quantities that we already encountered. Firstly,

$$
\begin{equation*}
K(\mathbf{X}) \lesssim \log \left(d_{\mathrm{BM}}\left(\ell_{2}^{n}, \mathbf{X}\right)+1\right) \lesssim \log n \tag{88}
\end{equation*}
$$

The first inequality in (88) is a useful theorem of Pisier [Pis80a, Pis80b]. The second inequality in (88) follows from John's theorem [Joh48], though for this purpose it suffices to use the older Auberbach lemma (see [Ban93, page 209] and Day47, Tay47]). By Pis80b] (see also e.g. [JS01, Lemma 17]) the rightmost quantity in (88) can be reduced if $\mathbf{X}$ is a subspace of $L_{1}$, namely we have

$$
\begin{equation*}
K(\mathbf{X}) \lesssim \mathrm{c}_{L_{1}}(\mathbf{X}) \sqrt{\log n} \tag{89}
\end{equation*}
$$

Secondly, $K(\mathbf{X})$ relates to the notion of type that we recalled in (78) through the following bounds:

$$
\begin{equation*}
T_{1+\frac{c}{K(\mathbf{X})^{2}}}(\mathbf{X})^{\frac{1}{2}} \lesssim K(\mathbf{X}) \leqslant \min _{p \in(1,2]} e^{\left(C T_{p}(\mathbf{X})\right)^{\frac{p}{p-1}}}, \tag{90}
\end{equation*}
$$

The qualitative meaning of (90) is that the $K$-convexity constant of a Banach space is finite if and only if it has type $p$ for some $p>1$; this is a landmark theorem of Pisier (the 'if' direction is due to [Pis82] and the 'only if' direction is due to [Pis73]). Since in our setting $\mathbf{X}$ is finite dimensional $(\operatorname{dim}(\mathbf{X})=n \geqslant 2)$, such a qualitative statement is vacuous without its quantitative counterpart (90). The first inequality in (90) can be deduced from [Pis83] (together with the computation of the implicit dependence on $p$ in Pis83] that was carried out in HLN16, Lemma 32]). The second inequality in (90) follows from an examination of the proof in [Pis82]. We omit the details of both deductions as they would result in a (quite lengthy and tedious) digression. It would be very interesting to determine the best bounds in the context of 90 .

Proposition 60 combined with 88 implies that Conjecture 10 holds up to a logarithmic factor in the sense that for every integer $n \geqslant 2$, any origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$ admits a matrix $S \in \mathrm{SL}_{n}(\mathbb{R})$ and an origin-symmetric convex body $L \subseteq S K$ such that

$$
\begin{equation*}
\frac{\operatorname{iq}(L)}{\sqrt{n}}\left(\frac{\operatorname{vol}_{n}(K)}{\operatorname{vol}_{n}(L)}\right)^{\frac{1}{n}} \lesssim \log n \tag{91}
\end{equation*}
$$

Furthermore, by (89) the $\log n$ in (91) can be replaced by $\sqrt{\log n}$ if $K$ is the unit ball of a subspace of $L_{1}$ (equivalently, the polar of $K$ is a zonoid), and by the second inequality in (90) if $p>1$, then the $\log n$ in (91) can be replaced by a dimension-independent quantity that depends only on $p$ and the type- $p$ constant of the norm whose unit ball is $K$. Also, Corollary 33 holds with the right hand side of (59) multiplied by $\log n$, and the reverse Faber-Krahn inequality of Conjecture 35 holds up to a factor of $(\log n)^{2}$, i.e., for any origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$ there is $S \in S L_{n}(\mathbb{R})$ such that $\lambda(S K) \operatorname{vol}(K)^{2 / n} \lesssim n(\log n)^{2}$. If $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a canonically positioned normed space, then it follows that for a suitable choice of normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ the right hand side of 28 , and hence also $e(\mathbf{X})$ by Theorem 21, is at most a universal constant multiple of $\operatorname{evr}(\mathbf{X}) \sqrt{n} \log n$, and also $n \lesssim \lambda(\mathbf{X}) \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{2 / n} \lesssim n(\log n)^{2}$.
1.7. Randomized clustering. All of the new upper bounds on Lipschitz extension moduli that we stated above rely on a geometric structural result for finite dimensional normed spaces (and subsets thereof). Beyond the application to Lipschitz extension, this result is of value in its own right because it yields an improvement of a basic randomized clustering method from the computer science literature.

The link between random partitions of metric spaces and Lipschitz extension was found in LN05. We will adapt the methodology of [LN05] to deduce the aforementioned Lipschitz extension theorems from our new bound on randomized partitions of normed spaces. In order to formulate the corresponding
definitions and results, one must first set some groundwork for a notion of a random partition of a metric space, whose subsequent applications necessitate certain measurability requirements.

A framework for reasoning about random partitions of metric spaces was developed in [LN05], but we will formulate a different approach. The reason for this is that the definitions of [LN05] are in essence the minimal requirements that allow one to use at once several different types of random partitions for Lipschitz extension, which leads to definitions that are more cumbersome than the approach that we take below. Greater simplicity is not the only reason why we chose to formulate a foundation that differs from [LN05]. The approach that we take is easier to implement, and, importantly, it yields a bi-Lipschitz invariant, while we do not know if the corresponding notions in [LN05] are bi-Lipschitz invariants (we suspect that they are not, but we did not attempt to construct examples that demonstrate this). The Lipschitz extension theorem of [LN05] is adapted accordingly in Section55 thus making the present article self-contained, and also yielding simplification and further applications. Nevertheless, the key geometric ideas that underly this use of random partitions are the same as in [LN05].

Obviously, there are no measurability issues when one considers finite metric spaces (in our setting, finite subsets of normed spaces). The ensuing measurability discussions can therefore be ignored in the finitary setting. In particular, the computer science literature on random partitions focuses exclusively on finite objects. So, for the purpose of algorithmic clustering, one does not need the more general treatment below, but it is needed for the purpose of Lipschitz extension.
1.7.1. Basic definitions related to random partitions. Let $\left(m, d_{m}\right)$ be a metric space. Suppose that $\mathcal{P} \subseteq 2^{m}$ is a partition of $m$. For $x \in M$, denote by $\mathcal{P}(x) \subseteq M$ the unique element of $\mathcal{P}$ to which $x$ belongs. The sets $\{\mathcal{P}(x)\}_{x \in m}$ are often called the clusters of $\mathcal{P}$. Given $\Delta>0$, one says that $\mathcal{P}$ is $\Delta$-bounded if diam $m(\mathcal{P}(x)) \leqslant \Delta$ for every $x \in M$, where $\operatorname{diam}_{m}(S)=\sup \left\{d_{m}(x, y): x, y \in S\right\}$ denotes the diameter of $\varnothing \neq S \subseteq m$.

Suppose that $(\mathscr{L}, \mathcal{F})$ is a measurable space, i.e., $\mathscr{L}$ is a set and $\mathcal{F} \subseteq 2^{\mathscr{Z}}$ is a $\sigma$-algebra of subsets of $\mathscr{L}$. Recall (see (Jac68] or the convenient survey |Wag77]) that if ( $m, d_{m}$ ) is a metric space, then a set-valued mapping $\Gamma: \mathscr{L} \rightarrow 2^{m}$ is said to be strongly measurable if for every closed subset $E \subseteq m_{\text {we have }}$

$$
\begin{equation*}
\Gamma^{-}(E) \stackrel{\text { def }}{=}\{z \in \mathscr{L}: E \cap \Gamma(z) \neq \varnothing\} \in \mathcal{F} \tag{92}
\end{equation*}
$$

Throughout what follows, when we say that $\mathcal{P}$ is a random partition of a metric space ( $m, d_{m}$ ), we mean the following (formally, the objects that we will be considering are random ordered partitions into countably many clusters). There is a probability space ( $\Omega, \mathbf{P r o b}$ ) and a sequence of set-valued mappings

$$
\left\{\Gamma^{k}: \Omega \rightarrow 2^{m}\right\}_{k=1}^{\infty}
$$

We write $\mathcal{P}^{\omega}=\left\{\Gamma^{k}(\omega)\right\}_{k=1}^{\infty}$ for each $\omega \in \Omega$ and require that the mapping $\omega \mapsto \mathcal{P}^{\omega}$ takes values in partitions of $m$. We also require that for every fixed $k \in \mathbb{N}$, the set-valued mapping $\Gamma^{k}: \Omega \rightarrow 2^{m}$ is strongly measurable, where the $\sigma$-algebra on $\Omega$ is the Prob-measurable sets. Given $\Delta>0$, we say that $\mathcal{P}$ is a $\Delta$-bounded random partition of $\left(m, d_{m}\right)$ if $\mathcal{P}^{\omega}$ is a $\Delta$-bounded partition of $\left(m, d_{m}\right)$ for every $\omega \in \Omega$.

Remark 61. Recall that when we say that $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a normed space we mean that the underlying vector space is $\mathbb{R}^{n}$, equipped with a norm $\|\cdot\|_{\mathrm{X}}: \mathbb{R}^{n} \rightarrow[0, \infty)$. By doing so, we introduce a second metric on $\mathbf{X}$, i.e., $\mathbb{R}^{n}$ is also endowed with the standard Euclidean structure that corresponds to the norm $\|\cdot\|_{\ell_{2}}$. This leads to ambiguity when we discuss $\Delta$-bounded partitions of $\mathbf{X}$ for some $\Delta>0$, as there are two possible metrics with respect to which one could bound the diameters of the clusters. In fact, a key aspect of our work is that it can be beneficial to consider another auxiliary norm $\|\cdot\|_{\mathrm{Y}}$ on $\mathbb{R}^{n}$, as in e.g. Theorem 21 , thus leading to three possible interpretations of $\Delta$-boundedness of a partition of $\mathbb{R}^{n}$. To avoid any confusion, we will adhere throughout to the convention that when we say that a partition $\mathcal{P}$ of $\mathbf{X}$ is $\Delta$-bounded we mean exclusively that all the clusters of $\mathcal{P}$ have diameter at most $\Delta$ with respect to the norm $\|\cdot\| \mathbf{x}$.
1.7.2. Iterative ball partitioning. Fix $\Delta \in(0, \infty)$. Iterative ball partitioning is a common procedure to construct a $\Delta$-bounded random partition of a metric probability space. We will next describe it to clarify at the outset the nature of the objects that we investigate, and because our new positive partitioning results
are solely about this type of partition. Thus, our contribution to the theory of random partitions is a sharp understanding of the performance of iterative ball partitioning of normed spaces, and, importantly, the demonstration of the utility of its implementation using balls that are induced by a suitably chosen auxiliary norm rather than the given norm that we aim to study. On the other hand, our impossibility results rule out the existence of any random partition whatsoever with certain desirable properties.

The iterative ball partitioning method is a ubiquitous tool in metric geometry and algorithm design. To the best of our knowledge, it was first used by Karger, Motwani and Sudan [KMS98] and the aforementioned work $\left[\widetilde{C C G}^{+} 98\right.$ in the context of normed spaces, and it has become very influential in the context of general metric spaces due to its use in that setting (with the important twist of randomizing the radii) by Calinescu, Karloff and Rabani [CKR05]. To describe it, suppose that ( $m, d_{m}$ ) is a metric space and that $\mu$ is a Borel probability measure on $m$. Let $\left\{\mathrm{X}_{k}\right\}$ inductively a sequence $\left\{\Gamma^{k}\right\}_{k=1}^{\infty}$ of random subsets of $m$ by setting $\Gamma^{1}=B_{m}\left(\mathrm{X}_{1}, \Delta / 2\right)$ and

$$
\forall k \in\{2,3, \ldots,\}, \quad \Gamma^{k} \stackrel{\text { def }}{=} B m\left(\mathrm{X}_{k}, \frac{\Delta}{2}\right) \backslash \bigcup_{j=1}^{k-1} B m\left(\mathrm{X}_{j}, \frac{\Delta}{2}\right) .
$$

By design, $\operatorname{diam} m\left(\Gamma^{k}\right) \leqslant \Delta$. Under mild assumptions on $m$ and $\mu$ that are simple to check, $\Gamma^{k}$ will have the measurability properties that we require below and $\mathcal{P}=\left\{\Gamma^{k}\right\}_{k=1}^{\infty}$ will be a partition of $M$ almost-surely. While initially the clusters of $\mathcal{P}$ are quite "tame," e.g. they start out as balls in $m$, as the iteration proceeds and we discard the balls that were used thus far, the resulting sets become increasingly "jagged." In particular, even when the underlying metric space ( $M, d_{m}$ ) is very "nice," the clusters of $\mathcal{P}$ need not be connected; see Figure 2. Nevertheless, we will see that such a simple procedure results in a random partition with probabilistically small boundaries in sense that will be described rigorously below.


FigURE 2. A schematic depiction of (randomized) iterative ball partitioning of a bounded subset of $\mathbb{R}^{2}$, where $\mathbb{R}^{2}$ is equipped with a norm whose unit ball is a regular hexagon. The centers of the above hexagons are chosen independently and uniformly at random from a large region that contains the given subset of $\mathbb{R}^{2}$. At each step of the iteration, a new hexagon appears, and it carves out a new cluster which consists of the part of the hexagon that does not intersect any of the clusters that have been formed in the previous stages of the iteration. The first few clusters that are formed by this procedure are typically hexagons, but at later stages the clusters become more complicated and less "round." In particular, they can eventually become disconnected, as exhibited by the region that is shaded black above.

In the present setting, the metric space that we wish to partition is a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$, so it is natural to want to use the Lebesgue measure on $\mathbb{R}^{n}$ in the above construction. Since this measure is not a probability measure, we cannot use the above framework directly. For this reason, we will in fact use a periodic variant of iterative ball partitioning of $\mathbf{X}$ by adapting a construction that was used in [LN05].
1.7.3. Separation and padding. Fix $\Delta>0$. Let $\mathcal{P}$ be a $\Delta$-bounded random partition of a metric space $M$. As a random "clustering" of $M$ into pieces of small diameter, $\mathcal{P}$ yields a certain "simplification" of $M$. For such a simplification to be useful, one must add a requirement that it "mimics" the geometry of $M$ in a meaningful way. The literature contains multiple definitions that achieve this goal, leading to applications in both algorithms and pure mathematics. We will not attempt to survey the literature on this topic, quoting only the definitions of separating and padded random partitions, which are the simplest and most popular notions of random partitions of metric spaces among those that have been introduced.

Definition 62 (separating random partition and separation modulus). Let $\left(M, d_{m}\right)$ be a metric space. For $\sigma, \Delta>0, a \Delta$-bounded random partition $\mathcal{P}$ of $\left(m, d_{m}\right)$ is $\sigma$-separating if

$$
\begin{equation*}
\forall x, y \in M, \quad \operatorname{Prob}[\mathcal{P}(x) \neq \mathcal{P}(y)] \leqslant \frac{\sigma}{\Delta} d_{m}(x, y) \tag{93}
\end{equation*}
$$

The separation modulu $\xi^{10}$ of $(m, d m)$, denoted $\operatorname{SEP}(m, d m)$ or simply $\mathrm{SEP}(M)$ if the metric is clear from the context, is the infimum over those $\sigma>0$ such that for every $\Delta>0$ there exists a $\sigma$-separating $\Delta$-bounded random partition of $(m, d m)$. If no such $\sigma$ exists, then write $\operatorname{SEP}(m, d m)=\infty$. Similarly, for $n \in \mathbb{N}$, the size-n separation modulus of $(m, d m)$, denoted $\operatorname{SEP}^{n}(m, d m)$ or simply $\operatorname{SEP}^{n}(m)$ if the metric is clear from the context, is the infimum over those $\sigma>0$ such that for every $S \subseteq M$ with $|S| \leqslant n$ and every $\Delta>0$ there exists a $\sigma$-separating $\Delta$-bounded random partition of $(S, d m)$. In other words,

$$
\operatorname{SEP}^{n}\left(m, d_{m}\right) \stackrel{\text { def }}{=} \sup _{\substack{S \subseteq m \\|S| \leqslant n}} \operatorname{SEP}\left(S, d_{m}\right)
$$

While the notions that we presented in Definition62 are standard (see below for the history), it will be beneficial for us (e.g. for proving Theorem 29) to introduce the following terminology.

Definition 63 (separation profile). Let $\left(M, d_{m}\right)$ be a metric space. We say that a metric $\mathfrak{d}: M \times M \rightarrow[0, \infty)$ on $M$ is a separation profile of $(m, d m)$ if for every $\Delta>0$ there exists a $\Delta$-bounded random partition $\mathcal{P}_{\Delta}$ of $\left(m, d_{m}\right)$ that is defined on some probability space $\left(\Omega_{\Delta}, \mathbf{P r o b}_{\Delta}\right)$ such that

$$
\begin{equation*}
\forall x, y \in M, \quad \mathfrak{d}(x, y) \geqslant \sup _{\Delta \in(0, \infty)} \Delta \operatorname{Prob}_{\Delta}\left[\mathcal{P}_{\Delta}(x) \neq \mathcal{P}_{\Delta}(y)\right] . \tag{94}
\end{equation*}
$$

So, the separation modulus of $(m, d m)$ is the infimum over those $\sigma>0$ for which $\sigma d m$ is a separation profile of $(M, d m)$. Definition 63 would make sense for functions $\mathfrak{d}: M \times M \rightarrow[0, \infty)$ that need not be metrics on $M$, but we prefer to deal only with separation profiles of $(m, d m)$ that are metrics on $M$ so as to be able to discuss the Lipschitz condition with respect to them; observe that the right hand side of (94) is a metric on $M$, so any such function is always at least (point-wise) a metric that is a separation profile of $\left(m, d_{m}\right)$. If $\mathfrak{d}: M \times M \rightarrow[0, \infty)$ is a separation profile of $\left(M, d_{m}\right)$, then $\mathfrak{d}(x, y) \geqslant d_{m}(x, y)$ for all $x, y \in M$ because $\operatorname{diam}_{m}\left(\mathcal{P}_{d_{m}(x, y)-\varepsilon}(x)\right) \leqslant d_{m}(x, y)-\varepsilon<d_{m}(x, y)$ for any $0<\varepsilon<d_{m}(x, y)$, so we necessarily have $y \notin \mathcal{P}_{d_{m}(x, y)-\varepsilon}(x)$ (deterministically) and therefore

$$
\begin{equation*}
\mathfrak{d}(x, y) \geqslant\left(d_{m}(x, y)-\varepsilon\right) \operatorname{Prob}_{d_{m}(x, y)-\varepsilon}\left[\mathcal{P}_{d_{m}(x, y)-\varepsilon}(x) \neq \mathcal{P}_{d_{m}(x, y)-\varepsilon}(y)\right]=d_{m}(x, y)-\varepsilon \tag{95}
\end{equation*}
$$

Definition 64 (padded random partition and padding modulus). Let ( $M, d_{m}$ ) be a metric space. For $\delta, \mathfrak{p}, \Delta>0, a \Delta$-bounded random partition $\mathcal{P}$ of $\left(M, d_{m}\right)$ is $(\mathfrak{p}, \delta)$-padded if

$$
\begin{equation*}
\forall x \in m, \quad \operatorname{Prob}\left[B m\left(x, \frac{\Delta}{\mathfrak{p}}\right) \subseteq \mathcal{P}(x)\right] \geqslant \delta \tag{96}
\end{equation*}
$$

[^6]Denote by $\mathrm{PAD}_{\delta}\left(m, d_{m}\right)$, or simply $\mathrm{PAD}_{\delta}(m)$ if the metric is clear from the context, the infimum over those $\mathfrak{p}>0$ such that for every $\Delta>0$ there exists $a(\mathfrak{p}, \delta)$-padded $\Delta$-bounded random partition $\mathcal{P}$ of $\left(m, d_{m}\right)$. If no such $\mathfrak{p}$ exists, then write $\operatorname{PAD}_{\delta}\left(m, d_{m}\right)=\infty$. For every $n \in \mathbb{N}$, denote

$$
\operatorname{PAD}_{\delta}^{n}\left(m, d_{m}\right) \stackrel{\text { def }}{=} \sup _{\substack{S \subseteq m \\|S| \leqslant n}} \operatorname{PAD}_{\delta}\left(S, d_{m}\right) .
$$

See Section 3 for a quick justification why the above definition of random partition implies that the events that appear in (93) and (96) are indeed Prob-measurable.

Qualitatively, condition (93) says that despite the fact that $\mathcal{P}$ decomposes $m$ into clusters of small diameter, any two nearby points are likely to belong to the same cluster. Condition (96) says that every point in $T$ is likely to be "well within" its cluster (its distance to the complement of its cluster is at least a definite proportion of the assumed upper bound on the diameter of that cluster). Both of these requirements express the (often nonintuitive) property that the "boundaries" that the random partition induces are "thin" in a certain distributional sense, despite the fact that each realization of the partition consists only of small diameter clusters that can sometimes be very jagged. Neither of the above two definitions implies the other, but it follows from [LN03] that if $\mathcal{P}$ is a $\mathfrak{p}, \delta)$-padded $\Delta$-bounded random partition of $\left(m, d_{m}\right)$, then there exits a random partition $\mathcal{P}^{\prime}$ of $\left(m, d_{m}\right)$ that is $(2 \Delta)$-bounded and ( $\left.4 \mathfrak{p} / \delta\right)$-separating.

Separating and padded random partitions were introduced in the articles Bar96, Bar99] of Bartal, which contained decisive algorithmic applications and influenced a flurry of subsequent works that obtained many more applications in several directions. Other works considered such partitions implicitly, with a variety of applications; see the works of Leighton-Rao [LR88], Awerbuch-Peleg [AP90], Linial-Saks [LS91], Alon-Karp-Peleg-West [AKPW91], Klein-Plotkin-Rao [KPR93] and Rao [Rao99]. The nomenclature of Definition 62 and Definition 64 comes from [GKL03, LN03, LN04a, LN05, KLMN05].

By Bar96], for every metric space $(m, d m)$ and every integer $n \geqslant 2$, we have $\operatorname{SEP}^{n}(m) \lesssim \log n$. It was observed by Gupta, Krauthgamer and Lee [GKL03] that [Bar96] also implicitly yields the padding bound $\mathrm{PAD}_{0.5}^{n}(m) \lesssim \log n$. It was proved in Bar96] that both of these estimates are sharp.

Random partitions of normed spaces were first studied by Peleg and Reshef [PR98] for applications to network routing and distributed computing. The aforementioned work $\left[\mathrm{CCG}^{+} 98\right]$ improved and generalized the bounds of [PR98], and influenced later works; see e.g. [LN05], and the work [A106] of Andoni and Indyk. Similar partitioning schemes appeared implicitly in earlier work [KMS98 on algorithms for graph colorings based on semidefinite programming.
1.7.4. From separation to Lipschitz extension. As we already explained, the connection between random partitions and Lipschitz extension was found in LN05. Here we will use the following theorem to deduce Theorem 29. It implies in particular the bound

$$
\begin{equation*}
\mathrm{e}(m) \lesssim \mathrm{SEP}(m) \tag{97}
\end{equation*}
$$

of [LN05] and its proof is an adaptation of the ideas of [LN05] to both the present setup (extension to a function that is Lipschitz with respect to a different metric) and our different measurability requirements from the random partitions; we stress, however, that even though we cannot apply [LN05 directly as a "black box," the geometric ideas that underly the proof of Theorem65 are the same as those of [LN05].

Theorem 65. Suppose that $\mathfrak{d}$ is a separation profile of a locally compact metric space ( $m, d_{m}$ ). For every Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ and every subset $\mathrm{C} \subseteq m$, iff $: \mathcal{C} \rightarrow \mathbf{Z}$ is 1 -Lipschitz with respect to the metric $d_{m}$, i.e., $\|f(x)-f(y)\|_{\mathbf{Z}} \leqslant d_{m}(x, y)$ for every $x, y \in m$, then there is $F: m \rightarrow \mathbf{Z}$ that extends $f$ and is $O(1)$-Lipschitz with respect to the metric $\mathfrak{d}$, i.e., $\|F(x)-F(y)\|_{\mathbf{z}} \lesssim \mathfrak{d}(x, y)$ for every $x, y \in m$.
1.7.5. Bounds on the separation and padding moduli of normed spaces. To facilitate the ensuing discussion of upper and lower bounds on the separation and padding moduli of (subsets of) normed spaces, we will first record two of their rudimentary properties. Firstly, the following lemma formally expresses
the aforementioned advantage of the definitions in Section 1.7.3 over those of [LN05], namely that the moduli $\operatorname{SEP}(\cdot)$ and $\operatorname{PAD}_{\delta}(\cdot)$ are bi-Lipschitz invariants; its straightforward proof appears in Section 3 .

Lemma 66 (bi-Lipschitz invariance of separation and padding moduli). Let ( $m, d_{m}$ ) be a complete metric space that admits a bi-Lipschitz embedding into a metric space $\left(n, d_{n}\right)$. Then

$$
\begin{equation*}
\operatorname{SEP}\left(m, d_{m}\right) \leqslant \mathrm{c}\left(n, d_{n}\right)\left(m, d_{m}\right) \operatorname{SEP}\left(\eta, d_{n}\right), \tag{98}
\end{equation*}
$$

and

$$
\begin{equation*}
\forall \delta \in(0,1), \quad \operatorname{PAD}_{\delta}\left(m, d_{m}\right) \leqslant \mathrm{c}_{\left(n, d_{n}\right)}\left(m, d_{m}\right) \mathrm{PAD}_{\delta}\left(n, d_{n}\right) . \tag{99}
\end{equation*}
$$

Secondly, we have the following tensorization property whose simple proof appears in Section3. For $s \in[1, \infty]$ and metric spaces $\left(m_{1}, d_{m_{1}}\right),\left(m, d_{m_{2}}\right)$, the metric $d_{m_{1} \oplus_{s}} m_{2}: m_{1} \times m_{2} \rightarrow[0, \infty)$ on the Cartesian product $m_{1} \times m_{2}$ is defined by setting for every $\left(x_{1}, x_{2}\right),\left(y_{1}, y_{2}\right) \in m_{1} \times m_{2}$,

$$
\begin{equation*}
d_{m_{1} \oplus_{s}} m_{2}\left(\left(x_{1}, x_{2}\right),\left(y_{1}, y_{2}\right)\right) \stackrel{\text { def }}{=}\left(d_{m}\left(x_{1}, y_{1}\right)^{s}+d_{n}\left(x_{2}, y_{2}\right)^{s}\right)^{\frac{1}{s}} . \tag{100}
\end{equation*}
$$

With the usual convention that when $s=\infty$ the right hand side of (100) is equal to the maximum of $d_{m}\left(x_{1}, y_{1}\right)$ and $d_{n}\left(x_{2}, y_{2}\right)$. The metric space $\left(m_{1} \times m_{2}, d_{m_{1} \oplus_{s} m_{2}}\right)$ is will be denoted $m_{1} \oplus_{s} m_{2}$.

Lemma 67 (tensorization of separation and padding moduli). For any $s \in[1, \infty]$ and $\delta_{1}, \delta_{2} \in(0,1)$, any two metric spaces $\left(m_{1}, d_{m_{1}}\right)$ and $\left(m_{2}, d_{m_{2}}\right)$ satisfy

$$
\begin{equation*}
\operatorname{SEP}\left(m_{1} \oplus_{s} m_{2}\right) \leqslant \operatorname{SEP}\left(m_{1}\right)+\operatorname{SEP}\left(m_{2}\right), \tag{101}
\end{equation*}
$$

and

$$
\begin{equation*}
\operatorname{PAD}_{\delta_{1} \delta_{2}}\left(m_{1} \oplus_{s} m_{2}\right) \leqslant\left(\operatorname{PAD}_{\delta_{1}}\left(m_{1}\right)^{s}+\operatorname{PAD}_{\delta_{2}}\left(m_{2}\right)^{s}\right)^{\frac{1}{s}} . \tag{102}
\end{equation*}
$$

The following theorem shows that the bi-Lipschitz invariant $\mathrm{PAD}_{\delta}(\cdot)$ is not sufficiently sensitive to distinguish substantially between normed spaces, as its value is essentially independent of the norm.

Theorem 68. For every $n \in \mathbb{N}$, every normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ satisfies

$$
\begin{equation*}
\forall \delta \in(0,1), \quad \frac{1}{1-\sqrt[n]{\delta}} \leqslant \frac{1}{2} \operatorname{PAD}_{\delta}(\mathbf{X}) \leqslant \frac{1+\sqrt[n]{\delta}}{1-\sqrt[n]{\delta}} . \tag{103}
\end{equation*}
$$

Therefore, $\operatorname{PAD}_{\delta}(\mathbf{X})=\max \left\{1, \frac{\operatorname{dim}(\mathbf{X})}{\log (1 / \delta)}\right\}$ for every finite dimensional normed space $\mathbf{X}$ and $\delta \in(0,1)$.
As we explained above, in the setting of Theorem 68 the fact that $\mathrm{PAD}_{0.5}(\mathbf{X})=O(n)$ is well-known. We will prove the upper bound on $\operatorname{PAD}_{\delta}(\mathbf{X})$ that appears in (103), i.e., with sharp dependence on both $n$ and $\delta$, in Section 4.1. The fact that $\mathrm{PAD}_{0.5}(\mathbf{X})$ is at least a universal constant multiple of $n$ was proved in the manuscript [LN03]. Because [LN03] is not intended for publication, we will prove the lower bound on $\mathrm{PAD}_{\delta}(\mathbf{X})$ that appears in (103) in Section 2.6. by following the reasoning of [N03] while taking more care than we did in [LN03] in order to obtain sharp dependence on $\delta$ in addition to sharp dependence on $n$.

In contrast to Theorem 68 the separation modulus of a finite dimensional normed space can have different asymptotic dependencies on its dimension. Indeed, $\operatorname{SEP}\left(\ell_{2}^{n}\right)=\sqrt{n}$ and $\operatorname{SEP}\left(\ell_{1}^{n}\right)=n$ by [CCG ${ }^{+} 98$ ], so using Lemma66 we see that every normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ satisfies the a priori bounds

$$
\begin{equation*}
\frac{n}{d_{\mathrm{BM}}\left(\ell_{1}^{n}, \mathbf{X}\right)} \lesssim \mathrm{SEP}(\mathbf{X}) \lesssim d_{\mathrm{BM}}\left(\ell_{2}^{n}, \mathbf{X}\right) \sqrt{n}, \tag{104}
\end{equation*}
$$

which we already quoted in the above overview as (2).
Giannopoulos proved Gia95] that every $n$-dimensional normed space $\mathbf{X}$ satisfies $d_{\mathrm{BM}}\left(\ell_{1}^{n}, \mathbf{X}\right) \lesssim n^{5 / 6}$, so the first inequality in (104) implies that $\operatorname{SEP}(\mathbf{X}) \gtrsim \sqrt[6]{n}$. Alternatively, the fact that $\operatorname{SEP}(\mathbf{X}) \geqslant n^{c}$ for some universal constant $c>0$ follows from by combining Theorem 1 with (97). Actually, we always have

$$
\begin{equation*}
\operatorname{SEP}(\mathbf{X}) \gtrsim \sqrt{n}, \tag{105}
\end{equation*}
$$

which coincides with the first half of (7). Observe that (105) cannot follow from a "vanilla" application of the first inequality in (104] by Szarek's work [Sza90]. In fact, the first inequality of (104) must sometimes yield a worse power type dependence on $n$ than in (105), because Tikhomirov proved in [Tik19] that there is a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ that satisfies $d_{\mathrm{BM}}\left(\ell_{1}^{n}, \mathbf{X}\right) \geqslant n^{a}$ for some universal constant $a>1 / 2$.

Nevertheless, we can prove (105) by the following a "hereditary" application of (104). Bourgain and Szarek [BS88] and independently Ball (see [BS88, Remark 7], [Sza91, Remark 7], [TJ89, page 138]) proved (relying on the Bourgain-Tzafriri restricted invertibility principle BT87]) that there is $m \in\{1, \ldots, n\}$ with $m=n$ such that $\mathrm{c}_{\mathbf{X}}\left(\ell_{1}^{m}\right) \lesssim \sqrt{n}$ (in fact, by BS88] any $2 n$-dimensional normed space has Banach-Mazur distance $O(\sqrt{n})$ from $\left.\ell_{1}^{n} \oplus \ell_{2}^{n}\right)$. Therefore, by (98) we have $\operatorname{SEP}(\mathbf{X}) \gtrsim \operatorname{SEP}\left(\ell_{1}^{m}\right) / c_{\mathbf{X}}\left(\ell_{1}^{m}\right)=m / c_{\mathbf{X}}\left(\ell_{1}^{m}\right) \gtrsim \sqrt{n}$.

The second half of (7) is the following lower bound on $\operatorname{SEP}(\mathbf{X})$ in terms of the type 2 constant of $\mathbf{X}$.

$$
\begin{equation*}
\operatorname{SEP}(\mathbf{X}) \gtrsim T_{2}(\mathbf{X})^{2} \tag{106}
\end{equation*}
$$

We will prove (106) in Section 2.2 using Talagrand's refinement [Tal92] of Elton's theorem [Elt83], by the same hereditary use of $[104)$, namely showing that there is $m \in\{1, \ldots, n\}$ for which $m / c_{\mathbf{X}}\left(\ell_{1}^{m}\right) \gtrsim T_{2}(\mathbf{X})^{2}$.

Remark 69. It is impossible to improve (7) for all the values of the relevant parameters, as seen by considering $\mathbf{X}=\ell_{2}^{n-m} \oplus_{2} \ell_{1}^{m}$ for each $m \in\{1, \ldots, n\}$. Indeed, since in this case $T_{2}(\mathbf{X})=\sqrt{m}$,

$$
\operatorname{SEP}(\mathbf{X}) \stackrel{101}{\leqslant} \operatorname{SEP}\left(\ell_{2}^{n-m}\right)+\operatorname{SEP}\left(\ell_{1}^{m}\right)=\sqrt{n-m}+m=\sqrt{n}+k=\max \left\{\sqrt{\operatorname{dim}(\mathbf{X})}, T_{2}(\mathbf{X})^{2}\right\}
$$

Thanks to [72, the following theorem is a restatement of the lower bound on $\operatorname{SEP}(\mathbf{X})$ in Theorem 3 ,
Theorem 70. For every $n \in \mathbb{N}$, any normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ satisfies $\operatorname{SEP}(\mathbf{X}) \gtrsim \operatorname{evr}(\mathbf{X}) \sqrt{n}$.
As $\operatorname{evr}(\mathbf{X}) \geqslant 1$ (by definition), Theorem 70 implies 105 , via a proof that differs from the above reasoning. Also, Theorem 70 is stronger than the first inequality in (104) because $\operatorname{evr}\left(\ell_{1}^{n}\right)=\sqrt{n}$, and hence

$$
\operatorname{evr}(\mathbf{X}) \sqrt{n} \geqslant \frac{\operatorname{evr}\left(\ell_{1}^{n}\right)}{d_{\mathrm{BM}}\left(\ell_{1}^{n}, \mathbf{X}\right)} \sqrt{n} \asymp \frac{n}{d_{\mathrm{BM}}\left(\ell_{1}^{n}, \mathbf{X}\right)}
$$

We will prove Theorem 70 in Section 2.5 by adapting to the setting of general normed spaces the strategy that was used in $\mathrm{CCG}^{+} 98$ to treat $\ell_{1}^{n}$. The volumetric lower bound on $\operatorname{SEP}(\mathbf{X})$ of Theorem 70 is typically quite easy to use and it often leads to estimates that are better than the first inequality in (104).

For example, by [Sch82, Proposition 2.2] the Schatten-von Neumann trace class $\mathrm{S}_{p}^{n}$ satisfies

$$
\begin{equation*}
\forall p \geqslant 1, \quad \operatorname{evr}\left(\mathrm{~S}_{p}^{n}\right)=n^{\max \left\{\frac{1}{p}-\frac{1}{2}, 0\right\}} \tag{107}
\end{equation*}
$$

By substituting (107) into Theorem 70 we get that

$$
\begin{equation*}
\forall 1 \leqslant p \leqslant 2, \quad \operatorname{SEP}\left(\mathrm{~S}_{p}^{n}\right) \gtrsim n^{\frac{1}{p}-\frac{1}{2}} \sqrt{\operatorname{dim}\left(\mathrm{~S}_{p}^{n}\right)}=n^{\frac{1}{p}+\frac{1}{2}} \tag{108}
\end{equation*}
$$

An upper bound that matches 108 is a consequence of the second inequality in 104 as follows

$$
\operatorname{SEP}\left(\mathrm{S}_{p}^{n}\right) \lesssim d_{\mathrm{BM}}\left(\mathrm{~S}_{p}^{n}, \ell_{2}^{n^{2}}\right) \sqrt{\operatorname{dim}\left(\mathrm{S}_{p}^{n}\right)}=d_{\mathrm{BM}}\left(\mathrm{~S}_{p}^{n}, \mathrm{~S}_{2}^{n}\right) n=n^{\frac{1}{p}+\frac{1}{2}}
$$

We therefore have

$$
\forall 1 \leqslant p \leqslant 2, \quad \operatorname{SEP}\left(\mathrm{~S}_{p}^{n}\right)=n^{\frac{1}{p}+\frac{1}{2}}
$$

At the same time, the first inequality in (104) does not imply (108) since by a theorem of Davis (which was published only in the monograph [TJ89]; see Theorem 41.10 there), for every $1 \leqslant p \leqslant 2$ we have

$$
\begin{equation*}
d_{\mathrm{BM}}\left(\ell_{1}^{n^{2}}, \mathrm{~S}_{p}^{n}\right)=n \tag{109}
\end{equation*}
$$

So, the first inequality in (104) only implies the weaker lower bound $\operatorname{SEP}\left(\mathrm{S}_{p}^{n}\right) \gtrsim n$. Of course, this rules out a "vanilla" use of $(104)$ and a hereditary application of $\sqrt[104]{ }$ as we did above could conceivably lead
to (108), i.e., there could be $m \in\{1, \ldots, n\}$ such that $m / \mathrm{c}_{p}^{n}\left(\ell_{1}^{m}\right)$ is at least the right hand side of 108). However, this possibility seems to be unlikely, as it would mean that the following conjecture has a negative answer, which would entail finding a remarkable (and likely valuable elsewhere) subspace of $\mathrm{S}_{p}^{n}$.

Conjecture 71. Fix $1 \leqslant p \leqslant 2$ and $0<\delta \leqslant 1$. If $n, m \in \mathbb{N}$ satisfy $m \geqslant \delta n^{2}$, then

$$
d_{\mathrm{BM}}\left(\ell_{1}^{m}, \mathbf{X}\right) \gtrsim_{p, \delta} n
$$

for every $m$-dimensional subspace $\mathbf{X}$ of $S_{p}^{n}$.
Thus, (109) is the case $\delta=1$ of Conjecture 71. which asserts that the same asymptotic lower bound persists if we consider subspaces of $S_{p}^{n}$ of proportional dimension rather than $S_{p}^{n}$ itself. Conjecture 71 is attractive in its own right, but it also implies that (108) does not follow from a hereditary application of the first inequality in (104). To see this, suppose for contradiction that there were $m \in\{1, \ldots, n\}$ such that

$$
\begin{equation*}
\frac{m}{\mathrm{CS}_{p}^{n}\left(\ell_{1}^{m}\right)} \gtrsim p n^{\frac{1}{p}+\frac{1}{2}} \tag{110}
\end{equation*}
$$

By Rademacher's differentiation theorem Rad19 there is an $m$-dimensional subspace $\mathbf{X}$ of $\mathrm{S}_{p}^{n}$ satisfying

$$
\begin{equation*}
\mathrm{c}_{p}^{n}\left(\ell_{1}^{m}\right)=d_{\mathrm{BM}}\left(\ell_{1}^{m}, \mathbf{X}\right) \gtrsim \frac{d_{\mathrm{BM}}\left(\ell_{1}^{m}, \ell_{2}^{m}\right)}{d_{\mathrm{BM}}\left(\mathrm{~S}_{p}^{n}, \mathrm{~S}_{2}^{n}\right)}=\frac{\sqrt{m}}{n^{\frac{1}{p}-\frac{1}{2}}} \tag{111}
\end{equation*}
$$

By contrasting (111) with (110) we deduce that necessarily $m \gtrsim p n^{2}$, so an application of Conjecture 71 gives $m / \mathrm{C}_{p}^{n}\left(\ell_{1}^{m}\right) \lesssim_{p} n$, which contradicts 110 since $p<2$.

Remark 72. The Löwner ellipsoid of $\ell_{\infty}^{n}\left(\ell_{1}^{n}\right)$ is $\sqrt{n} B_{\ell_{2}^{n}\left(\ell_{2}^{n}\right)}$, and $B_{\ell_{\infty}^{n}\left(\ell_{1}^{n}\right)}=\left(B_{\ell_{1}^{n}}\right)^{n}$. Consequently,

$$
\operatorname{evr}\left(\ell_{\infty}^{n}\left(\ell_{1}^{n}\right)\right) n=n\left(\frac{(\pi n)^{\frac{n^{2}}{2}} / \Gamma\left(\frac{n^{2}}{2}+1\right)}{2^{n^{2}} /(n!)^{n}}\right)^{\frac{1}{n^{2}}}=n^{\frac{3}{2}}
$$

Therefore, Theorem 70 gives

$$
\begin{equation*}
\operatorname{SEP}\left(\ell_{\infty}^{n}\left(\ell_{1}^{n}\right)\right) \gtrsim n^{\frac{3}{2}} \tag{112}
\end{equation*}
$$

We will soon see that (112) is optimal, though unlike the above discussion for $\mathrm{S}_{p}^{n}$ when $1 \leqslant p \leqslant 2$, this does not follow from the second inequality in (104) because by (KS89],

$$
\begin{equation*}
d_{\mathrm{BM}}\left(\ell_{2}^{n^{2}}, \ell_{\infty}^{n}\left(\ell_{1}^{n}\right)\right)=d_{\mathrm{BM}}\left(\ell_{1}^{n^{2}}, \ell_{\infty}^{n}\left(\ell_{1}^{n}\right)\right)=n \tag{113}
\end{equation*}
$$

(113) also shows that (112) does not follow from the first inequality in (104). It seems that the method used in [KS89] to prove (113) is insufficient for proving that (112) does not follow from a hereditary application of the first inequality in (104). Analogously to Conjecture 71, we conjecture that this is impossible, which is a classical-sounding question about Banach-Mazur distances of independent interest.

Before passing to a description of our upper bounds on the separation modulus, we formulate the following corollary of Theorem 70 on the separation modulus of norms whose unit ball is a polytope; it restates the lower bound (6) and establishes its optimality.

Theorem 73. Fix $n \in \mathbb{N}$ and a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$. Suppose that $B_{\mathbf{X}}$ is a polytope that has exactly $\rho n$ vertices (note that necessarily $\rho \geqslant 2$, since $B_{\mathbf{X}}$ is origin-symmetric). Then

$$
\begin{equation*}
\operatorname{SEP}(\mathbf{X}) \gtrsim \frac{n}{\sqrt{\log \rho}} \tag{114}
\end{equation*}
$$

Moreover, this bound cannot be improved in general.

As an example of a consequence of Theorem 73] let $\mathbf{G}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{G}}\right)$ be a Gluskin space [Glu81], i.e. it is a certain random norm on $\mathbb{R}^{n}$ whose unit ball has $O(n)$ vertices; see the survey [MTJ03] for extensive information about this important construction and its variants. The expected Banach-Mazur distance between two independent copies of $\mathbf{G}$ is at least $c n$ for some universal constant $c>0$, so the expected Banach-Mazur distance between $\mathbf{G}$ and $\ell_{1}^{n}$ is at least $\sqrt{c n}$. Thus, the first inequality in (104) only shows that $\operatorname{SEP}(\mathbf{G}) \gtrsim \sqrt{n}$ in expectation, while Theorem 73 shows that in fact $\operatorname{SEP}(\mathbf{G}) \gtrsim n / \sqrt{\log n}$. It would be interesting to determine the growth rate of $\mathbb{E}[\operatorname{SEP}(\mathbf{G})]$. In particular, can it be that $\mathbb{E}[\operatorname{SEP}(\mathbf{G})] \gtrsim n$ ?

Proof of Theorem 73. By applying a linear isometry of $\mathbf{X}$ we may assume that $B_{\ell_{2}^{n}}$ is the Löwner ellipsoid of $B_{\mathbf{X}}$. Since $B_{\mathbf{X}}$ is a polytope with $\rho n$ vertices that is contained in $B_{\ell_{2}^{n}}$, we have

$$
\sqrt[n]{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} \lesssim \frac{\sqrt{\log \rho}}{n}
$$

by a result of Maurey [Pis81] (see also [Car85, BF87, CP88, Glu88, BLM89, BP90, Kyr00] and the expository treatments in Bal01, BGVV14]). Hence, evr(X) $\gtrsim \sqrt{n / \log \rho}$, so (114) follows from Theorem 70 .

Consider the following (dual of an) example of Figiel and Johnson [FJ80]. Fix $m \in \mathbb{N}$. Let $\mathbf{Z}=\left(\mathbb{R}^{m},\|\cdot\|_{\mathbf{z}}\right)$ be a normed space with $d_{\mathrm{BM}}\left(\ell_{2}^{m}, \mathbf{Z}\right) \lesssim 1$ such that $B_{\mathbf{Z}}$ is a polytope of $e^{O(m)}$ vertices; e.g. $B_{\mathbf{Z}}$ can be taken to be the convex hull of a net of $S^{m-1}$. For $k \in \mathbb{N}$, let $\mathbf{X}=\ell_{1}^{k}(\mathbf{Z})$. So, $\operatorname{dim}(\mathbf{X})=k m$ and $B_{\mathbf{X}}$ is a polytope of $2 k e^{O(m)}$ vertices. Thus (114) becomes $\operatorname{SEP}(\mathbf{X}) \gtrsim k \sqrt{m}$. At the same time, since $d_{\mathrm{BM}}\left(\ell_{2}^{m}, \mathbf{Z}\right) \lesssim 1$ we have $d_{\mathrm{BM}}\left(\ell_{2}^{k m}, \mathbf{X}\right) \lesssim \sqrt{k}$, so by (104) in fact $\operatorname{SEP}(\mathbf{X}) \lesssim \sqrt{k} \cdot \sqrt{k m}=k \sqrt{m}$, i.e., (114) is sharp in this case.

Theorem 29 follows from Theorem 65 thanks to the following randomized partitioning theorem.
Theorem 74. For every $n \in \mathbb{N}$ and every normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$, the metric $\mathfrak{d}$ that is defined by

$$
\forall x, y \in \mathbb{R}^{n}, \quad \mathfrak{d}(x, y)=\frac{4\|x-y\|_{\Pi^{*} \mathbf{x}}}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}
$$

is a separation profile for $\mathbf{X}$.
To illustrate Theorem 74 , fix $1 \leqslant p \leqslant \infty$ and apply it when $\mathbf{X}$ is the space $\mathbf{Y}_{p}^{n}$ of Theorem 24. By using Theorem 74 we see that for every $\Delta>0$ there is a random partition $\mathcal{P}$ of $\mathbb{R}^{n}$ with the following properties.
(1) For every $x \in \mathbb{R}^{n}$ we have $\operatorname{diam}_{\ell_{p}^{n}}(\mathcal{P}(x)) \leqslant \Delta$.
(2) For every $x, y \in \mathbb{R}^{n}$ we have

$$
\begin{equation*}
\operatorname{Prob}[\mathcal{P}(x) \neq \mathcal{P}(y)] \lesssim \frac{\|x-y\|_{\Pi^{*}} \mathbf{Y}_{p}^{n}}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}_{p}^{n}}\right.} \stackrel{300}{\lesssim} \stackrel{399}{\lesssim} \frac{n^{\frac{1}{p}}}{\Delta}\|x-y\|_{\ell_{2}^{n}} . \tag{115}
\end{equation*}
$$

In comparison to the $O(\sqrt{n})$-separating partition of $\ell_{2}^{n}$ from [CCG 98 ], when $p<2$ the above random partition has smaller clusters in the sense that their diameter in the $\ell_{p}^{n}$ metric is at most $\Delta$, which is more stringent than the requirement that their Euclidean diameter is at most $\Delta$. This improved control on the size of the clusters comes at the cost that in the probabilistic separation requirement (115) the quantity that multiplies the Euclidean distance increases from $O(\sqrt{n})$ to $O\left(n^{1 / p}\right)$. When $p>2$ this tradeoff is reversed, i.e., we get an asymptotic improvement in the separation guarantee (115) at the cost of requiring less from the cluster size, namely the diameter of each cluster is now guaranteed to be small in the $\ell_{p}^{n}$ metric rather than the more stringent requirement that it is small in the Euclidean metric.

Theorem 75 below follows from Theorem 74 the same way we deduced Theorem 21 from Theorem 29 ,
Theorem 75. Fix $n \in \mathbb{N}$ and two normed spaces $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right), \mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$. Every closed $\mathrm{C} \subseteq \mathbb{R}^{n}$ satisfies

$$
\begin{equation*}
\operatorname{SEP}\left(\mathcal{C}_{\mathbf{X}}\right) \leqslant 4\left(\sup _{\substack{x, y \in \mathbb{C} \\ x \neq y}} \frac{\|x-y\|_{\mathbf{X}}}{\|x-y\|_{\mathbf{Y}}}\right) \sup _{\substack{x, y \in \mathcal{C} \\ x \neq y}}\left(\frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{(x-y)^{\perp}}\left(B_{\mathbf{Y}}\right)\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \cdot \frac{\|x-y\|_{\ell_{2}^{n}}}{\|x-y\|_{\mathbf{X}}}\right) . \tag{116}
\end{equation*}
$$

Proof of Theorem 75 assuming Theorem 74. Let $M, M^{\prime}$ be as in (51). By Theorem 74 applied to $\mathbf{Y}$, for every $\Delta>0$ there is a random partition $\mathcal{P}$ of $\mathbb{R}^{n}$ that is $(\Delta / M)$-bounded with respect to $\mathbf{Y}$, i.e.,

$$
\frac{\operatorname{diam}_{\mathbf{X}}(\mathcal{P}(x))}{M} \stackrel{51}{\leqslant} \operatorname{diam}_{\mathbf{Y}}(\mathcal{P}(x)) \leqslant \frac{\Delta}{M}
$$

for every $x \in \mathbb{R}^{n}$, and also, recalling Definition 63. for every distinct $x, y \in \mathbb{R}^{n}$ we have

The special case $\mathcal{C}=\mathbb{R}^{n}$ of Theorem 75 coincides (with an explicitly stated constant factor) with the upper bound on $\operatorname{SEP}(\mathbf{X})$ in Theorem 3 since under the normalization $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$ we have

$$
\operatorname{SEP}(\mathbf{X}) \stackrel{\sqrt[30]{30} \wedge \sqrt{116}}{\leqslant} 4 \frac{\sup _{z \in \partial B_{\mathbf{X}}}\|z\|_{\Pi^{*} \mathbf{Y}}}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)}=4 \frac{\left\|\operatorname{ld}_{n}\right\|_{\mathbf{X} \rightarrow \Pi^{*} \mathbf{Y}}}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)}=4 \frac{\left\|\operatorname{ld}_{n}\right\|_{\Pi \mathbf{Y} \rightarrow \mathbf{X}^{*}}}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)}=2 \frac{\operatorname{diam}_{\mathbf{X}^{*}}\left(\Pi B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} .
$$

Also, Theorem 75 is stronger than the second inequality in (104) because by applying a linear isometry of $\mathbf{X}$ we may assume without loss of generality that $\|x\|_{\mathbf{X}} \leqslant\|x\|_{\ell_{2}^{n}} \leqslant d_{\mathrm{BM}}\left(\ell_{2}^{n}, \mathbf{X}\right)\|x\|_{\mathbf{X}}$ for all $x \in \mathbb{R}^{n}$, in which case the special case $\mathcal{C}=\mathbb{R}^{n}$ and $\mathbf{Y}=\ell_{2}^{n}$ of (116) implies that

$$
\operatorname{SEP}(\mathbf{X}) \leqslant \frac{4 \operatorname{vol}_{n-1}\left(B_{\ell_{2}^{n-1}}\right)}{\operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right)} d_{\mathrm{BM}}\left(\ell_{2}^{n}, \mathbf{X}\right)=\frac{4 \pi^{\frac{n-1}{2}} \Gamma\left(\frac{n}{2}+1\right)}{\pi^{\frac{n}{2}} \Gamma\left(\frac{n-1}{2}+1\right)} d_{\mathrm{BM}}\left(\ell_{2}^{n}, \mathbf{X}\right)=\frac{2^{\frac{3}{2}}+o(1)}{\sqrt{\pi}} d_{\mathrm{BM}}\left(\ell_{2}^{n}, \mathbf{X}\right) \sqrt{n} .
$$

The right hand side of (116) coincides (up to a universal constant factor) with the right hand side of (28), so all of the upper bounds for the Lipschitz extension modulus that we derived in the previous sections from Theorem 21 hold for the separation modulus, by Theorem 75 . For the separation modulus, we get several lower bounds from Theorem 70 that either provably match our upper bounds up to lower order factors, or match them assuming our conjectural isomorphic reverse isoperimetry. We will next spell out some of those consequences on randomized clustering of high dimensional norms.

Theorem 76. For every $p \geqslant 1, n \in \mathbb{N}$ and $k, r \in\{1, \ldots, n\}$ we have

$$
\begin{equation*}
\operatorname{SEP}\left(\left(\ell_{p}^{n}\right)_{\leqslant k}\right)=k^{\max \left\{\frac{1}{p}, \frac{1}{2}\right\}}, \tag{117}
\end{equation*}
$$

and

$$
r^{\max \left\{\frac{1}{p}, \frac{1}{2}\right\}} \sqrt{n} \lesssim \operatorname{SEP}\left(\left(\mathrm{~S}_{p}^{n}\right) \leqslant r\right) \lesssim r^{\max \left\{\frac{1}{p}, \frac{1}{2}\right\}} \sqrt{n} \cdot \begin{cases}\sqrt{\max \left\{\log \left(\frac{n}{r}\right), p\right\}} & \text { if } p \leqslant \log r  \tag{118}\\ \sqrt{\log n} & \text { if } p \geqslant \log r\end{cases}
$$

Moreover, if Conjecture 48 holds for $\mathbf{X}=\mathrm{S}_{p}^{n}$, then in fact

$$
\begin{equation*}
\operatorname{SEP}\left(\left(S_{p}^{n}\right) \leqslant r\right)=r^{\max \left\{\frac{1}{p}, \frac{1}{2}\right\}} \sqrt{n} . \tag{119}
\end{equation*}
$$

Proof. The deduction of the upper bounds on the separation modulus in (117) and (118) from Theorem 75 are identical, respectively, to the ways we deduced Theorem 20 and (45) from Theorem 21 .

For the first inequality in (117), since $\left(\ell_{p}^{n}\right) \leqslant k$ contains an isometric copy of $\ell_{p}^{k}$, we have

$$
\operatorname{SEP}\left(\left(\ell_{p}^{n}\right)_{\leqslant k}\right) \geqslant \operatorname{SEP}\left(\ell_{p}^{k}\right) \gtrsim \frac{k}{d_{\mathrm{BM}}\left(\ell_{p}^{k}, \ell_{1}^{k}\right)} \stackrel{104}{\frac{104}{2}} \frac{k}{k^{\max \left\{1-\frac{1}{p}, \frac{1}{2}\right\}}}=k^{\min \left\{\frac{1}{p}, \frac{1}{2}\right\}},
$$

where the asymptotic evaluation of $d_{\mathrm{BM}}\left(\ell_{p}^{k}, \ell_{q}^{k}\right)$ for all $p, q \geqslant 1$ is due Gurariĭ, Kadec' and Macaev [GKM66].
For the first inequality in [118], use the fact that $\left(\mathrm{S}_{p}^{n}\right)_{\leqslant r}$ contains an isometric copy of $\mathrm{S}_{p}^{r \times n}$, which is the Schatten-von Neumann trace class on the $r$-by- $n$ real matrices $\mathrm{M}_{r \times n}(\mathbb{R})$, whose norm is given by

$$
\begin{equation*}
\forall A \in M_{r \times n}(\mathbb{R}), \quad\|A\|_{S_{p}^{r \times n}}=\left(\operatorname{Tr}\left(\left(A A^{*}\right)^{\frac{p}{2}}\right)\right)^{\frac{1}{p}} . \tag{120}
\end{equation*}
$$

We then have the following rectangular version of (107) whose derivation is explained in Remark 170 .

$$
\begin{equation*}
\operatorname{evr}\left(\mathrm{S}_{p}^{r \times n}\right)=r^{\max \left\{\frac{1}{p}-\frac{1}{2}, 0\right\}} . \tag{121}
\end{equation*}
$$

The desired lower bound on $\operatorname{SEP}\left(\left(S_{p}^{n}\right)_{\leqslant r}\right)$ is now an application of Theorem 70 .
Remark 77. Theorem 3.3 in $\left[\mathrm{CCG}^{+} 98\right]$ asserts that $\operatorname{SEP}\left(\ell_{p}^{n}\right)=n^{\max \{1 / p, 1-1 / p\}}$ for every $p \geqslant 1$. If $p \geqslant 2$, then this means that it was previously thought that $\operatorname{SEP}\left(\ell_{p}^{n}\right)=n^{1-1 / p}$, which is of course incompatible with the case $k=n$ of (117) for every $p>2$. While $\mathrm{CCG}^{+} 98$ provides a complete proof of $\operatorname{SEP}\left(\ell_{p}^{n}\right)=n^{1 / p}$ when $1 \leqslant p \leqslant 2$, in the range $p>2$ the assertion $\operatorname{SEP}\left(\ell_{p}^{n}\right)=n^{1-1 / p}$ in [CCG 98 is justified through the use of a result from reference [14] in [CCG ${ }^{+} 98$, which is cited there as a "personal communication" with P. Indyk (dated April 1998). This reference was never published; after discovering the improved estimate of Theorem76, we confirmed with Indyk that it was indeed flawed.

Corollary 78. Conjecture 48 implies Conjecture 6. Namely, if Conjecture 48 holds for a canonically positioned normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$, then

$$
\begin{equation*}
\operatorname{SEP}(\mathbf{X})=\operatorname{evr}(\mathbf{X}) \sqrt{n} \simeq \operatorname{vr}\left(\mathbf{X}^{*}\right) \sqrt{n} . \tag{122}
\end{equation*}
$$

In particular, if $\mathbf{X}$ satisfies the assumptions of Lemma 52 (e.g. if $\mathbf{X}$ is symmetric), then (122) holds. Furthermore, if $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ is a symmetric normed space, then $\operatorname{SEP}\left(\mathrm{S}_{\mathbf{E}}\right)=\operatorname{evr}(\mathbf{E}) n^{1+o(1)}$. More precisely,

$$
\begin{equation*}
\operatorname{evr}(\mathbf{E}) n \lesssim \operatorname{SEP}\left(\mathrm{~S}_{\mathbf{E}}\right) \lesssim \operatorname{evr}(\mathbf{E}) n \sqrt{\log n} . \tag{123}
\end{equation*}
$$

Proof. The lower bound on $\operatorname{SEP}(\mathbf{X})$ in (122) is Theorem 70 (thus, it requires neither Conjecture 48 nor $\mathbf{X}$ being canonically positioned). The matching upper bound on $\operatorname{SEP}(\mathbf{X})$ in (122) follows from Corollary 50 and the fact that by Theorem 75 the separation modulus of any (not necessarily canonically positioned) normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is bounded from above by the right hand side of (53). The rest of the assertions of Corollary 78 follow from Lemma 52 and Lemma 53 .

By incorporating Proposition 60 into the same reasoning as in the justification of Corollary 78 , we also deduce the following stronger version of Theorem 12 .
Theorem 79. If $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a canonically positioned normed space, then

$$
\operatorname{evr}(\mathbf{X}) \sqrt{n} \lesssim \operatorname{SEP}(\mathbf{X}) \lesssim K(\mathbf{X}) \operatorname{evr}(\mathbf{X}) \sqrt{n} \stackrel{88}{\lesssim} \operatorname{evr}(\mathbf{X}) \sqrt{n} \log n .
$$

Section 6.3 contains volume ratio computations that show how Corollary 78 and Theorem 79 imply Corollary 4 as well as the conjectural (i.e., conditional on the validity of Conjecture 48 for the respective spaces) asymptotic evaluations (14) and (15), and several further results of this type. Most of the volume ratio computations in Section 6.3 rely on the available literature (notably Schütt's work [Sch82]), with a few new twists that are perhaps of independent geometric/probabilisitic interest (e.g. Lemma 172).
1.7.6. Dimension reduction. Fix $n \in \mathbb{N}$ and a metric space $\left(m, d_{m}\right)$. Recall that in Definition 62 we denoted by $\operatorname{SEP}^{n}\left(m, d_{m}\right)$ the supremum over all the separation moduli of subsets of $m$ of size at most $n$. In $\left[\mathrm{CCG}^{+} 98\right]$ it was shown that $\mathrm{SEP}^{n}\left(\ell_{2}\right) \lesssim \sqrt{\log n}$. Indeed, this follows from the Johnson-Lindenstrauss dimension reduction lemma [JL84], which asserts that any $n$-point subset of $\ell_{2}$ can be embedded with $O(1)$ distortion into $\ell_{2}^{m}$ with $m \lesssim \log n$, combined with the proof in $\mathrm{CCG}^{+} 98$ that $\operatorname{SEP}\left(\ell_{2}^{m}\right) \lesssim \sqrt{m}$.

One might expect that the optimal bounds that we know for $\operatorname{SEP}\left(\ell_{p}^{n}\right)$ in the entire range $p \in(1, \infty)$ also translate to improved bounds on $\operatorname{SEP}^{n}\left(\ell_{p}\right)$. The term "improved" is used here to mean any upper bound of the form $o_{p}(\log n)$ as $n \rightarrow \infty$, since the benchmark general result is the aforementioned upper bound $\operatorname{SEP}^{n}\left(m, d_{m}\right) \lesssim \log n$ from [Bar96], which holds for any $n$-point metric space $\left(m, d_{m}\right)$. This bound is sharp in general Bar96, so (because every $n$-point metric space embeds isometrically into $\ell_{\infty}^{n}$ ) we cannot hope to get a better bound on $\operatorname{SEP}^{n}\left(\ell_{\infty}\right)$ despite the fact that we obtained here an improved upper bound on $\operatorname{SEP}\left(\ell_{\infty}^{n}\right)$.

The obstacle is that when $p \in[1, \infty] \backslash\{2\}$ no bi-Lipschitz dimension reduction result is known for finite subsets of $\ell_{p}$, and poly-logarithmic bi-Lipschitz dimension reduction is impossible if $p \in\{1, \infty\}$; the case $p=\infty$ is due to Matoušek [Mat96] (see also [Nao17b, Nao21a]) and the case $p=1$ is due to Brinkman and Charikar [BC05] (see also [LN04b, Reg13, NPS20, NY21]). When $p \in[1, \infty] \backslash\{1,2, \infty\}$ remarkably nothing is known, i.e., neither positive results nor impossibility results are available for bi-Lipschitz dimension reduction, and it is a major open problem to make any progress in this setting; see [Naol8] for more on this area. Despite this obstacle, we have the following theorem that treats the range $p \in[1,2]$.

Theorem 80. For every $p \in(1,2]$ and $n \in \mathbb{N}$ we have

$$
(\log n)^{\frac{1}{p}} \lesssim \operatorname{SEP}^{n}\left(\ell_{p}\right) \lesssim \frac{(\log n)^{\frac{1}{p}}}{p-1}
$$

The lower bound on $\operatorname{SEP}^{n}\left(\ell_{p}\right)$ of Theorem 80 can be deduced from [CCG $\left.{ }^{+} 98\right]$; see Section 4.2 for the details. An upper bound of $\operatorname{SEP}^{n}\left(\ell_{p}\right) \lesssim_{p}(\log n)^{T / p}$ was obtained when $p \in(1,2]$ in the manuscript [LN03]. As [LN03] is not intended for publication, a proof of the upper bound on $\operatorname{SEP}^{n}\left(\ell_{p}\right)$ that is stated in Theorem 80 is included in Section 4.2 where we perform the argument with more care than the way we initially did it in [N03], so as to obtain the best dependence on $p$ that is achievable by this approach. Nevertheless, we conjecture that the dependence on $p$ in Theorem 80 could be removed altogether, though this would likely require a substantially new idea.

Conjecture 81. The dependence on $p$ in Theorem 80 can be improved to $\operatorname{SEP}^{n}\left(\ell_{p}\right) \lesssim(\log n)^{\frac{1}{p}}$.
So, if $p \leqslant 1+c(\log \log \log n) / \log \log n$ for some universal constant $c>0$, then Theorem 80 does not improve asymptotically over $\operatorname{SEP}^{n}\left(\ell_{p}\right) \lesssim \log n$, while Conjecture 81 would imply that $\operatorname{SEP}^{n}\left(\ell_{p}\right)=o(\log n)$ if and only if $\lim _{n \rightarrow \infty}(p-1) \log \log n=\infty$.

For fixed $p \in(2, \infty)$, at present we do not see how to obtain an upper bound on $\operatorname{SEP}^{n}\left(\ell_{p}\right)$ of the form $o_{p}(\log n)$ as $n \rightarrow \infty$. We state this separately as an interesting and challenging open question.

Question 82. Is it true that for every $n \in \mathbb{N}$ and $p \in(2, \infty)$ we have $\lim _{n \rightarrow \infty} \operatorname{SEP}^{n}\left(\ell_{p}\right) / \log n=0$ ? More ambitiously, is it true that $\operatorname{SEP}^{n}\left(\ell_{p}\right) \lesssim p \sqrt{\log n}$ ?

Note that $\operatorname{SEP}^{n}(\mathbf{X}) \gtrsim \sqrt{\log n}$ for any infinite-dimensional normed space $\mathbf{X}$, because by Dvoretzky's theorem Dvo61] we have $\mathrm{c}_{\mathbf{X}}\left(\ell_{2}^{m}\right)=1$ for every $m \in \mathbb{N}$, and therefore $\operatorname{SEP}^{n}(\mathbf{X}) \geqslant \operatorname{SEP}^{n}\left(\ell_{2}\right)=\sqrt{\log n}$.
1.8. Consequences in the linear theory. Even though the purpose of the present article was to investigate the nonlinear invariants e( $\cdot$ ) and $\operatorname{SEP}(\cdot)$, by relating them to volumetric quantities and other linear invariants of Banach spaces (such as type and cotype), we arrive at consequences that have nothing to do with nonlinear issues. In this section, we will give a flavor of such consequences, though we will not be exhaustive since it would be more natural to pursue them separately for their own right in future work.

Denote the Minkowski functional of an origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$ by $\|\cdot\|_{K}$, i.e., it is the norm on $\mathbb{R}^{n}$ whose unit ball is equal to $K$. The following theorem coincides with the second inequality in (1) upon a straightforward application of duality as we did in (31); this formulation is intended to highlight how we are bounding a convex-geometric quantity by a bi-Lipschitz invariant.

Theorem 83 (nonsandwiching between a convex body and its polar projection body). Fix $n \in \mathbb{N}$ and $\alpha, \beta \in(0, \infty)$. Let $K, L \subseteq \mathbb{R}^{n}$ be symmetric convex bodies with $\operatorname{vol}_{n}(L)=1$. Suppose that

$$
\begin{equation*}
\alpha L \subseteq K \subseteq \beta \Pi^{*} L . \tag{124}
\end{equation*}
$$

Then,

$$
\begin{equation*}
\frac{\beta}{\alpha} \gtrsim \operatorname{SEP}\left(\mathbb{R}^{n},\|\cdot\|_{K}\right) . \tag{125}
\end{equation*}
$$

Since the separation modulus of a metric space is at least the separation modulus of any of its subsets, by combining (125) with the first inequality in (1) we see that the sandwiching hypothesis (124) implies the following purely volumetric consequence for every linear subspace $\mathbf{V} \subseteq \mathbb{R}^{n}$.

$$
\begin{equation*}
\frac{\beta}{\alpha} \gtrsim \operatorname{evr}(K \cap \mathbf{V}) \sqrt{n}=\operatorname{vr}\left(\operatorname{Proj}_{\mathbf{V}} K^{\circ}\right) \sqrt{n} . \tag{126}
\end{equation*}
$$

In particular, using $\operatorname{evr}\left(\ell_{1}^{n}\right) \asymp \sqrt{n}$, we record separately the following special case of 126.
Corollary 84 (nonsandwiching of the cross-polytope). Fix $n \in \mathbb{N}$ and $\alpha, \beta \in(0, \infty)$. If $L \subseteq \mathbb{R}^{n}$ is a convex body of volume 1 that satisfies $\alpha L \subseteq B_{\ell_{1}^{n}} \subseteq \beta \Pi^{*} L$, then necessarily $\beta / \alpha \gtrsim n$.

The geometric meaning of Theorem 83 when $L=K$ is spelled out in the following corollary.
Corollary 85 (every origin-symmetric convex body admits a large cone). For every $n \in \mathbb{N}$, every originsymmetric convex body $K \subseteq \mathbb{R}^{n}$ has a boundary point $z \in \partial K$ that satisfies

$$
\begin{equation*}
\frac{\operatorname{vol}_{n}\left(\operatorname{Cone}_{z}(K)\right)}{\operatorname{vol}_{n}(K)} \gtrsim \frac{1}{n} \operatorname{SEP}\left(\mathbb{R}^{n},\|\cdot\|_{K}\right) . \tag{127}
\end{equation*}
$$

To see that Corollary 85 coincides with the case $L=K$ of Theorem 83 , simply recall the definition of the polar projection body $\Pi^{*} K$ in (30), while also recalling that for $z \in \mathbb{R}^{n} \backslash\{0\}$ we denote the cone whose base is $\operatorname{Proj}_{z^{\perp}}(K) \subseteq z^{\perp}$ and whose apex is $z$ by $\operatorname{Cone}_{z}(K)$, and the volume of Cone $(K)$ is given in (35).

A substitution of (105) into Corollary 85 shows that any origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$ has a boundary point $z \in \partial K$ that satisfies

$$
\begin{equation*}
\frac{\operatorname{vol}_{n}\left(\operatorname{Cone}_{z}(K)\right)}{\operatorname{vol}_{n}(K)} \gtrsim \frac{1}{\sqrt{n}} \tag{128}
\end{equation*}
$$

It seems (based on inquiring with experts in convex geometry) that the classical-looking geometric statement (128) did not previously appear in the literature. However, in response to our inquiry Lutwak found a different proof of $(128)$ which in addition shows that the best possible constant in $(128)$ is $1 / \sqrt{2 \pi}$. More precisely, we have the following proposition, whose proof (which relies on classical Brunn-Minkowski theory, unlike the indirect way by which we found (128)), is included in Section 2.7 (this proof is a restructuring of the proof that Lutwak found; we thank him for allowing us to include it here).

Proposition 86 (Lutwak). For every $n \in \mathbb{N}$, any origin symmetric convex body $K \subseteq \mathbb{R}^{n}$ satisfies

$$
\begin{equation*}
\max _{z \in \partial K} \frac{\operatorname{vol}_{n}\left(\operatorname{Cone}_{z}(K)\right)}{\operatorname{vol}_{n}(K)} \geqslant \frac{\Gamma\left(\frac{n}{2}\right)}{2 \sqrt{\pi} \Gamma\left(\frac{n+1}{2}\right)} \geqslant \frac{1+\frac{1}{4 n}}{\sqrt{2 \pi n}} . \tag{129}
\end{equation*}
$$

Moreover, the first inequality in (129) holds as equality if and only if K is an ellipsoid.
A substitution of (106) into Corollary 85yields the following geometric inequality.
Corollary 87. Fix $n \in \mathbb{N}$ and suppose that $K \subseteq \mathbb{R}^{n}$ is an origin-symmetric convex body. There is a boundary point $z \in \partial K$ such that the following inequality holds for every $x_{1}, \ldots, x_{n} \in K$.

$$
\begin{equation*}
\frac{\operatorname{vol}_{n}\left(\operatorname{Cone}_{z}(K)\right)}{\operatorname{vol}_{n}(K)} \gtrsim \frac{1}{n} f_{S^{n-1}}\left\|\sum_{i=1}^{n} \theta_{i} x_{i}\right\|_{K}^{2} \mathrm{~d} \theta \tag{130}
\end{equation*}
$$

By combining [TJ79] with Lemma 101 below, the maximum of the right hand side of (130) over all possible $x_{1}, \ldots, x_{n} \in K$ is bounded above and below by universal constant multiples of $T_{2}\left(\mathbb{R}^{n},\|\cdot\|_{K}\right)^{2} / n$ (recall the definition (78) of the type-2 constant), so Corollary 87 is indeed a substitution of (106) into 127).

Returning to Corollary 85 recall that both the cross-polytope $B_{\ell_{1}^{n}}$ and the hypercube $[-1,1]^{n}$ are examples of extremal symmetric convex bodies $K \subseteq \mathbb{R}^{n}$ that have a boundary point $z \in \partial K$ for which the volume of $\operatorname{Cone}_{z}(K)$ is a universal constant proportion of the volume of $K$ (the Euclidean ball is an example of a convex body that is not extremal in this regard). But, there is a difference between the cross-polytope and the hypercube in terms of the stability of this property. Specifically, there is an origin-symmetric
convex body $K \subseteq[-1,1]^{n} \subseteq O(1) K$ such that for every $z \in \partial K$ the left hand side of (127) is at most a universal constant multiple of $1 / \sqrt{n}$. In contrast, the following proposition shows that the extremality of $\max _{z \in \partial B_{\ell_{1}^{n}}} \operatorname{vol}_{n}\left(\operatorname{Cone}_{z}\left(B_{\ell_{1}^{n}}\right)\right) / \operatorname{vol}_{n}\left(B_{\ell_{1}^{n}}\right)$ (up to constant factors) persists under $O(1)$ perturbations.
Proposition 88. Fix $n \in \mathbb{N}$ and $\alpha, \beta \in(0, \infty)$. Suppose that $K \subseteq \mathbb{R}^{n}$ is an origin-symmetric convex body that satisfies $\alpha K \subseteq B_{\ell_{1}^{n}} \subseteq \beta K$. Then there exists a boundary point $z \in \partial K$ such that

$$
\begin{equation*}
\frac{\operatorname{vol}_{n}\left(\operatorname{Cone}_{z}(K)\right)}{\operatorname{vol}_{n}(K)} \gtrsim \frac{\alpha}{\beta} . \tag{131}
\end{equation*}
$$

Proposition 88 is a direct consequence of Corollary 85 , the bi-Lipschitz invariance of the modulus of separated decomposability, and the lower bound $\operatorname{SEP}\left(\ell_{1}^{n}\right) \gtrsim n$ of [CCG 98 .

The following proposition is an application in a different direction of the results that we described in the preceding sections.

Proposition 89. If $\left(\mathbf{E},\|\cdot\|_{\mathbf{E}}\right)$ is a finite dimensional normed space with a 1 -symmetric basis, then every subspace $\mathbf{X}$ of $\mathbf{E}$ satisfies

$$
\begin{equation*}
\operatorname{evr}(\mathbf{X}) \sqrt{\operatorname{dim}(\mathbf{X})} \lesssim \operatorname{evr}(\mathbf{E}) \sqrt{\operatorname{dim}(\mathbf{E})} . \tag{132}
\end{equation*}
$$

Proposition 89 holds because $\operatorname{SEP}(\mathbf{E}) \lesssim \operatorname{evr}(\mathbf{E}) \sqrt{\operatorname{dim}(\mathbf{E})}$ by Corollary 78 , while $\operatorname{SEP}(\mathbf{X}) \gtrsim \operatorname{evr}(\mathbf{X}) \sqrt{\operatorname{dim}(\mathbf{X})}$ by Theorem 70 so sollows from $\operatorname{SEP}(\mathbf{X}) \leqslant \operatorname{SEP}(\mathbf{E})$. This justification shows that Proposition 89 holds for a class of spaces that is larger than those that have a 1 -symmetric basis, and Conjecture 6 would imply that Proposition 89 holds when $\mathbf{E}$ is any canonically positioned normed space.

Nevertheless, Proposition 89 fails to hold without any further assumption on the normed space E. For example, the computation in Remark 51 shows that for any $n, m \in \mathbb{N}$ with $n \geqslant 2$ and $m=n \log n$, the space $\mathbf{E}=\ell_{1}^{n} \oplus \ell_{2}^{m}$ satisfies $\operatorname{evr}(\mathbf{E}) \sqrt{\operatorname{dim}(\mathbf{E})} \lesssim \sqrt{n \log n}$ while its subspace $\mathbf{X}=\ell_{1}^{n}$ satisfies $\operatorname{evr}(\mathbf{X}) \sqrt{\operatorname{dim}(\mathbf{X})}=n$.

Proposition 89 shows that if $\mathbf{E}$ has a 1 -symmetric basis, then among the linear subspaces $\mathbf{X}$ of $\mathbf{E}$ the invariant $\operatorname{evr}(\mathbf{X}) \sqrt{\operatorname{dim}(\mathbf{X})}$ is maximized up to universal constant factors at $\mathbf{X}=\mathbf{E}$. The fact we are multiplying here the external volume ratio of $\mathbf{X}$ by the square root of its dimension is an artifact of our proof and it would be interesting to understand what correction factors allow for such a result to hold:

Question 90. Characterize (up to universal constant factors) those $A:[1, \infty) \rightarrow[1, \infty)$ with the property that for any $n \geqslant 1$ we have $\operatorname{evr}(\mathbf{X}) A(k) \leqslant \operatorname{evr}(\mathbf{E}) A(n)$ for every normed space $\left(\mathbf{E},\|\cdot\|_{\mathbf{E}}\right)$ of dimension at most $n$ that has a 1 -symmetric basis, every $k \in\{1, \ldots, n\}$, and every $k$-dimensional subspace $\mathbf{X}$ of $\mathbf{E}$.

Proposition 89 shows that if $A(n)=\sqrt{n}$, then $A:[1, \infty) \rightarrow[1, \infty)$ has the properties that are described in Question 90, At the same time, no $A:[1, \infty) \rightarrow[1, \infty)$ with $A(n)=O(1)$ can be as in Question 90. Indeed, for any such $A$ consider the symmetric normed space $\mathbf{E}=\ell_{\infty}^{n}$. There is a universal constant $\eta>0$ such that any normed space $\mathbf{X}$ with $\operatorname{dim}(\mathbf{X}) \leqslant \eta \log n$ is at Banach-Mazur distance at most 2 from a subspace of $\ell_{\infty}^{n}{ }^{11}$ In particular, this holds for $\mathbf{X}=\ell_{1}^{m}$ when $m \in \mathbb{N}$ satisfies $m \leqslant \eta \log n$, so we get that

$$
\begin{equation*}
A(\eta \log n) \sqrt{\log n}=\operatorname{evr}\left(\ell_{1}^{m}\right) A(\eta \log n) \leqslant 2 \operatorname{evr}\left(\ell_{\infty}^{n}\right) A(n)=A(n) \tag{133}
\end{equation*}
$$

So, $A(n) \gtrsim \sqrt{\log n}$ and by iterating (133) one gets the slightly better lower bound $A(n) \gtrsim \sqrt{(\log n) \log \log n}$, as well as $A(n) \gtrsim \sqrt{(\log n)(\log \log n) \log \log \log n}$ and so forth, yielding in the end the estimate

$$
\begin{equation*}
A(n) \geqslant \frac{\left(\prod_{k=1}^{\log ^{*} n} \log ^{[k]} n\right)^{\frac{1}{2}}}{e^{O\left(\log ^{*} n\right)}} \tag{134}
\end{equation*}
$$

[^7]where for $k \in \mathbb{N} \cup\{0\}$ we denote the $k^{\prime}$ th iterant of the logarithm by $\log ^{[k]}$, i.e., $\log ^{[0]} x=x$ for $x>0$, and
\[

$$
\begin{equation*}
\log ^{[k]} x>0 \Longrightarrow \log ^{[k+1]} x=\log \left(\log ^{[k]} x\right) . \tag{135}
\end{equation*}
$$

\]

There is no reason to expect that the lower bound (134) is close to being optimal, but in combination with Proposition 89it does show that the answer to Question 90 is likely nontrivial.

These considerations lead to the following open-ended question. The literature contains results showing that $\ell_{p}^{n}$ maximizes certain geometric invariants (e.g. Banach-Mazur distance to $\ell_{2}^{n}$ [Lew78], or volume ratio (Bal91c]) among all the $n$-dimensional subspaces or quotients of $L_{p}$. Is there an analogous theory in the spirit of $(132)$ in the much more general setting of spaces that have a 1 -symmetric basis? This could be viewed as a symmetric space variant of the classical work of Lewis [Lew78, Lew79]. An interesting step in this direction can be found in [TJ80]; specifically, see [T]80] Theorem 1.2], which could be relevant to Question 90 through the approach of [Bal91c, Section 2].

Acknowledgements. I am grateful to Moses Charikar and Piotr Indyk for a helpful conversation on the erroneous optimality assertion of $\left[\mathrm{CCG}^{+} 98\right]$ in the range $p \in(2, \infty]$, and to Piotr Indyk for subsequent conversations and pointers to the literature. I thank Gideon Schechtman for multiple discussions that led in particular to Lemma 172 and (86). I thank Emmanuel Breuillard for discussions regarding Remark 40 and for sharing with me his proofs of statements on the comparison between the notions of being canonically positioned and having enough symmetries. I thank Erwin Lutwak for showing me his proof of Proposition 86, I thank Bo'az Klartag and Emanuel Milman for showing me the proof of Proposition 60, I am grateful to Keith Ball, Franck Barthe, Ronen Eldan, Charles Fefferman, Apostolos Giannopoulos, David Jerison, Grigoris Paouris, Gilles Pisier, Oded Regev, Carsten Schütt, Ramon van Handel and Elisabeth Werner for helpful discussions and suggestions.

## 2. Lower bounds

In this section we will prove the impossibility results that were stated in the Introduction. Throughout what follows, all Banach spaces will be tacitly assumed to be separable. Given a Banach space $\mathbf{X}$, its Banach-Mazur distance to a Hilbert space will be denoted $\mathrm{d}_{\mathbf{X}} \in[1, \infty]$, i.e., $\mathrm{d}_{\mathbf{X}}=d_{\mathrm{BM}}(\mathbf{X}, \mathbf{H})$ where $\mathbf{H}$ is a Hilbert space with either $\operatorname{dim}(\mathbf{H})=\operatorname{dim}(\mathbf{X})$ when $\operatorname{dim}(\mathbf{X})<\infty$, or $\mathbf{H}=\ell_{2}$ when $\mathbf{X}$ is infinite dimensional. By a classical result of Enflo [Enf70, Theorem 6.3.3] (see also [BL00, Corollary 7.10]) we have $\mathrm{d}_{\mathbf{X}}=\mathrm{c}_{2}(\mathbf{X})$.
2.1. Proof of Theorem 13. Recall that the (Gaussian) type 2 and cotype 2 constants of a Banach space $(\mathbf{X},\|\cdot\| \mathbf{x})$, denoted $T_{2}(\mathbf{X})$ and $C_{2}(\mathbf{X})$, respectively, are the infimum over those $T \in[1, \infty]$ and $C \in[1, \infty]$, respectively, for which the following inequalities hold for every $m \in \mathbb{N}$ and every $x_{1}, \ldots, x_{m} \in \mathbf{X}$ :

$$
\begin{equation*}
\frac{1}{C^{2}} \sum_{j=1}^{m}\left\|x_{j}\right\|_{\mathbf{X}}^{2} \leqslant \mathbb{E}\left[\left\|\sum_{j=1}^{m} \mathrm{~g}_{j} x_{j}\right\|_{\mathbf{X}}^{2}\right] \leqslant T^{2} \sum_{j=1}^{m}\left\|x_{j}\right\|_{\mathbf{X}}^{2} \tag{136}
\end{equation*}
$$

where henceforth $g_{1}, g_{2}, \ldots$ will always denote i.i.d. standard Gaussian random variables. The following theorem of Kwapień 【Kwa72〕 is fundamental (see also [Pis86, Theorem 3.3] or [TJ89, Theorem 13.15]).

Theorem 91. Every Banach space $\left(\mathbf{X},\|\cdot\|_{\mathbf{x}}\right)$ satisfies $\mathrm{d}_{\mathbf{X}} \leqslant T_{2}(\mathbf{X}) C_{2}(\mathbf{X})$.
We will use Theorem 91 to estimate the following quantity, which in turn will be used to get the best bound that we currently have on the constant $c$ that appears in the lower bound on $\mathrm{e}(\mathbf{X})$ of Theorem 13 ,

Definition 92 (Lindenstrauss-Tzafriri constant). Suppose that $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ is a Banach space. Define LT $(\mathbf{X})$ to be the infimum over those $K \in[1, \infty]$ such that for every closed linear subspace $\mathbf{V} \subseteq \mathbf{X}$ there exists a projection $\operatorname{Proj}: \mathbf{X} \rightarrow \mathbf{V}$ from $\mathbf{X}$ onto $\mathbf{V}$ whose operator norm satisfies $\|\operatorname{Proj}\|_{\mathbf{X} \rightarrow \mathbf{x}} \leqslant K$.

So, the Lindenstrauss-Tzafriri constant of a Hilbert space equals 1, and Sobczyk proved [Sob41] that

$$
\begin{equation*}
\forall n \in \mathbb{N}, \quad \operatorname{LT}\left(\ell_{1}^{n}\right)=\mathrm{LT}\left(\ell_{\infty}^{n}\right)=\sqrt{n} . \tag{137}
\end{equation*}
$$

We chose the nomenclature of Definition 92 in reference to the famous solution [LT71] by Lindenstrauss and Tzafriri of the complemented subspace problem, which asserts that if $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ is a Banach space for which $L T(\mathbf{X})<\infty$, then $\mathbf{X}$ is isomorphic to a Hilbert space, i.e., $\mathrm{d}_{\mathbf{X}}<\infty$. Moreover, if $\mathbf{X}$ is infinite dimensional, then it was shown in [LT71] that $\mathrm{d}_{\mathbf{X}} \lesssim \mathrm{LT}(\mathbf{X})^{4}$. This dependence was improved in [KM73] by Kadec and Mitjagin, who established the following theorem, which is the currently best-known bound in the Lindenstrauss-Tzafriri theorem (see also [Fig77, Pis88, Pis96, AK06, Kal08] for subsequent improvements of the implicit universal constant factor and further generalizations).

Theorem 93. Every infinite dimensional Banach space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ satisfies $\mathrm{d}_{\mathbf{X}} \lesssim \mathrm{LT}(\mathbf{X})^{2}$.
When $\operatorname{dim}(\mathbf{X})<\infty$ the question of bounding $d_{\mathbf{X}}$ by a function of $\operatorname{LT}(\mathbf{X})$ was left open in [LT71]. This question, which was eventually solved by Figiel, Lindenstrauss and Milman [FLM77, Theorem 6.7], turned out to be significantly more subtle than its infinite dimensional counterpart. The currently best-known estimate is due to Tomczak-Jaegermann [TJ89, Theorem 29.4], who proved the following theorem.
Theorem 94. Every finite dimensional Banach space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ satisfies $\mathrm{d}_{\mathbf{X}} \lesssim \mathrm{LT}(\mathbf{X})^{5}$.
The proof of Theorem 94 is achieved in [TJ89] through an interesting combination of the proof of the Lindenstrauss-Tzafriri theorem [LT71] with the finite dimensional machinery of [FLM77] and Milman's Quotient of Subspace Theorem Mil85.

The following theorem is a link between the Lindenstrauss-Tzafriri constant and Lipschitz extension.
Theorem 95. Every Banach space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ satisfies $\mathrm{e}(\mathbf{X}) \geqslant \mathrm{LT}(\mathbf{X})$.
Proof. By Remark 97 , if $\operatorname{dim}(\mathbf{X})=\infty$, then $\mathrm{e}(\mathbf{X})=\infty$, so we may assume that $\operatorname{dim}(\mathbf{X})<\infty$. Fix $L>\mathrm{e}(\mathbf{X})$ and let $\mathbf{V} \subseteq \mathbf{X}$ be a linear subspace of $\mathbf{X}$. Then, the identity mapping from $\mathbf{V}$ to $\mathbf{V}$ can be extended to an $L$-Lipschitz mapping $\rho: \mathbf{X} \rightarrow \mathbf{V}$. In other words, $\rho$ is an $L$-Lipschitz retraction from $\mathbf{X}$ onto $\mathbf{V}$. By a classical theorem of Lindenstrauss Lin64 (see also its elegant alternative proof by Pełczyńsky in Peł68, page 61]), there is a projection of norm at most $L$ from $\mathbf{X}$ onto $\mathbf{V}$. This proves that $\mathrm{LT}(\mathbf{X}) \leqslant L$.

The following theorem is the lower bound $\mathrm{e}\left(\ell_{2}^{n}\right) \gtrsim \sqrt[4]{n}$ of [MN13] that we already quoted in 22], in combination with the bi-Lipschitz invariance of the Lipschitz extension modulus.
Theorem 96. For every $n \in \mathbb{N}$, any normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ satisfies $\mathrm{e}(\mathbf{X}) \gtrsim \frac{\sqrt[4]{n}}{d_{\mathbf{x}}}$.
Remark 97. The question whether $\mathrm{e}\left(\ell_{2}\right)$ is finite or infinite was open for quite some time: It was first stated in print in [JLS86, page 137], and it was also posed by Ball in Bal92, page 170] (Ball conjectured that $\left.\mathrm{e}\left(\ell_{2}\right)=\infty\right)$. We answered it in Nao01] by proving that $\lim _{n \rightarrow \infty} \mathrm{e}\left(\ell_{2}^{n}\right)=\infty$. Due to Dvoretzky's theorem [Dvo61] this implies that $e(\mathbf{X})$ is at least an unbounded function of $\operatorname{dim}(\mathbf{X})$ for any normed space $\mathbf{X}$, and in particular $\mathrm{e}(\mathbf{X})=\infty$ if $\operatorname{dim}(\mathbf{X})=\infty$. A rate at which $\mathrm{e}\left(\ell_{2}^{n}\right)$ tends to $\infty$ was not specified in [Nao01], but the reasoning of Nao01] was inspected quantitatively in LN05, Remark 5.3], yielding an explicit lower bound that depends on an auxiliary parameter, and it was noted in $\overline{\mathrm{BB} 07 \mathrm{a}}$ that an optimization over this parameter yields the estimate $\mathrm{e}\left(\ell_{2}^{n}\right) \gtrsim \sqrt[8]{n}$. A further improvement from MN13] (whose proof refines ideas of Kalton Kal04, Kal12]) was the aforementioned estimate e $\left(\ell_{2}^{n}\right) \gtrsim \sqrt[4]{n}$ (a different proof of this bound follows from [Nao21b]), which is the currently best-known lower bound on e( $\ell_{2}^{n}$ ). By Milman's sharpening Mil71 of Dvoretzky's theorem Dvo61], it follows that every normed space $\mathbf{X}$ satisfies $\mathrm{e}(\mathbf{X}) \gtrsim \sqrt[4]{\log n}$. As we explained in Section 1.3 , the bound $\mathrm{e}\left(\ell_{\infty}^{n}\right) \gtrsim \sqrt{n}$ is classical (specifically, by substituting (137) into Theorem 95. In combination with the Alon-Milman theorem AM83] (see also [Tal95]), the fact that both $\mathrm{e}\left(\ell_{2}^{n}\right)=n^{\Omega(1)}$ and $\mathrm{e}\left(\ell_{\infty}^{n}\right)=n^{\Omega(1)}$ formally implies that

$$
\mathrm{e}(\mathbf{X}) \geqslant e^{\eta \sqrt{\log n}}
$$

for some universal constant $\eta>0$ and every $n$-dimensional normed space $\mathbf{X}$, which was the best-known general lower bound on the Lipschitz extension modulus prior to Theorem 1 .

The above results imply as follows the lower bound on $e(\mathbf{X})$ of Theorem 13. By combining Theorem 94 and Theorem 95 , we have $e(\mathbf{X}) \gtrsim \sqrt[5]{d_{\mathbf{X}}}$. In combination with Theorem 96 it therefore follows that

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \gtrsim \max \left\{\frac{\sqrt[4]{n}}{\mathrm{~d}_{\mathbf{X}}}, \sqrt[5]{\mathrm{d}_{\mathbf{X}}}\right\} \geqslant \sqrt[24]{n}, \tag{138}
\end{equation*}
$$

where the last step follows from elementary calculus and holds as equality when $\mathrm{d}_{\mathbf{X}}=n^{5 / 24}$.
We will derive a better lower bound on $\mathrm{e}(\mathbf{X})$ than $(138$ through the following theorem which improves over the power of $\mathrm{LT}(\mathbf{X})$ in Theorem 94 showing that in the finite dimensional setting one can come close (up to logarithmic factors) to the infinite dimensional bound of Theorem 93; see also Remark 102 below.
Theorem 98. For every integer $n \geqslant 2$, any $n$-dimensional Banach space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ satisfies

$$
\begin{equation*}
\mathrm{d}_{\mathbf{X}} \lesssim \mathrm{LT}(\mathbf{X})^{2}(\log n)^{3} . \tag{139}
\end{equation*}
$$

Assuming Theorem [138, reason analogously to (138) while using (139) in place of Theorem 94 to get

$$
\begin{equation*}
\mathrm{e}(\mathbf{X}) \gtrsim \max \left\{\frac{\sqrt[4]{n}}{\mathrm{~d}_{\mathbf{X}}}, \frac{\sqrt{\mathrm{d}_{\mathbf{X}}}}{(\log n)^{3}}\right\} \geqslant \frac{n^{\frac{1}{12}}}{(\log n)^{2}}, \tag{140}
\end{equation*}
$$

where equality holds in the final step of $(140)$ if and only if $\mathrm{d} \mathbf{x}=\sqrt[6]{n}(\log n)^{2}$.
Prior to proving Theorem 98 , we will record the following two standard lemmas that will be used in its proof; both will be established in correct generality that also treats infinite dimensional Banach spaces even though here we will need them only in the finite dimensional setting (the infinite dimensional formulations are relevant to the discussion in Remark 102).
Lemma 99. For every Banach space $\left(\mathbf{X},\|\cdot\|_{\mathbf{x}}\right)$ we have $\mathrm{LT}\left(\mathbf{X}^{*}\right) \leqslant \mathrm{LT}(\mathbf{X})+1$.
Proof. We may assume that $\mathrm{LT}(\mathbf{X})<\infty$. Then $\mathbf{X}$ is reflexive (even isomorphic to Hilbert space), by [LT71]. Fix a closed linear subspace $\mathbf{W}$ of $\mathbf{X}^{*}$ and denote its pre-annihilator by

$$
{ }^{\perp} \mathbf{W} \stackrel{\text { def }}{=} \bigcap_{x^{*} \in \mathbf{W}}\left\{x \in \mathbf{X}: x^{*}(x)=0\right\} \subseteq \mathbf{X} .
$$

Suppose that $K>\operatorname{LT}(\mathbf{X})$. By the definition of $\operatorname{LT}(\mathbf{X})$ there exists $\operatorname{Proj}: \mathbf{X} \rightarrow \mathbf{X}$ that is a projection from $\mathbf{X}$ onto ${ }^{\perp} \mathbf{W}$ whose operator norm satisfies $\|\operatorname{Proj}\|_{\mathbf{X} \rightarrow \mathbf{X}} \leqslant K$. Observe that for every $x^{*} \in \mathbf{X}^{*}$ and $x \in{ }^{\perp} \mathbf{W}$,

$$
\left(x^{*}-\operatorname{Proj}^{*} x^{*}\right)(x)=x^{*}(x)-x^{*}(\operatorname{Proj} x)=0,
$$

since $\operatorname{Proj} x=x$. This shows that

$$
\left(\mathrm{Id}_{\mathbf{X}^{*}}-\operatorname{Proj}^{*}\right)\left(\mathbf{X}^{*}\right) \subseteq\left({ }^{\perp} \mathbf{W}\right)^{\perp}=\left\{x^{*} \in X^{*}: x^{*}\left({ }^{\perp} \mathbf{W}\right)=\{0\}\right\}=\mathbf{W},
$$

where the last step follows from the double annihilator theorem since $\mathbf{X}$ is reflexive and hence $\mathbf{W}$ is weak* closed in $\mathbf{X}^{*}$. If $x^{*} \in \mathbf{W}$, then for any $x \in \mathbf{X}$ we have $\operatorname{Proj}^{*} x^{*}(x)=x^{*}(\operatorname{Proj} x)=0$, as $\operatorname{Proj} x \in{ }^{\perp} \mathbf{W}$. Hence $\operatorname{Proj}^{*} x^{*}=0$, and so $\operatorname{ld}_{\mathbf{X}^{*}}-$ Proj$^{*}$ acts as the identity when it is restricted to $\mathbf{W}$, i.e., $\operatorname{ld}_{\mathbf{X}^{*}}-\operatorname{Proj}^{*}: \mathbf{X}^{*} \rightarrow \mathbf{X}^{*}$ is a projection from $\mathbf{X}^{*}$ onto $\mathbf{W}$. It remains to note that

$$
\left\|\mathrm{Id}_{\mathbf{X}^{*}}-\operatorname{Proj}^{*}\right\|_{\mathbf{X}^{*} \rightarrow \mathbf{X}^{*}} \leqslant 1+\left\|\operatorname{Proj}^{*}\right\|_{\mathbf{X}^{*} \rightarrow \mathbf{X}^{*}}=1+\|\operatorname{Proj}\|_{\mathbf{X} \rightarrow \mathbf{X}} \leqslant K+1 .
$$

The following simple lemma shows that the Lindenstrauss-Tzafriri constant is a bi-Lipschitz invariant.
Lemma 100. Any two Banach spaces $\left(\mathbf{W},\|\cdot\|_{\mathbf{w}}\right)$ and $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ satisfy

$$
\begin{equation*}
\operatorname{LT}(\mathbf{W}) \leqslant c_{\mathbf{X}}(\mathbf{W}) \operatorname{LT}(\mathbf{X}) . \tag{141}
\end{equation*}
$$

Proof. We may assume that $\mathrm{c}_{\mathbf{X}}(\mathbf{W})<\infty$ and $\mathrm{LT}(\mathbf{X})<\infty$. By [LT7], the latter assumption implies that $\mathbf{X}$ is isomorphic to a Hilbert space, and hence it is reflexive. We may therefore apply a differentiation argument (see e.g. [BL00, Corollary 7.10]) to deduce that there is a closed subspace $\mathbf{Y}$ of $\mathbf{X}$ such that $d_{\mathrm{BM}}(\mathbf{W}, \mathbf{Y})=\mathrm{c}_{\mathbf{X}}(\mathbf{W})$. In other words, for every $D>\mathrm{c}_{\mathbf{X}}(\mathbf{W})$ there is a linear isomorphism $T: \mathbf{W} \rightarrow \mathbf{Y}$ satisfying $\|T\|_{\mathrm{W} \rightarrow \mathbf{Y}}\left\|T^{-1}\right\|_{\mathbf{Y} \rightarrow \mathbf{W}}<D$. If $\mathbf{V}$ is a closed subspace of $\mathbf{W}$ and $K>\mathrm{LT}(\mathbf{X})$, then there is a projection Proj from $\mathbf{X}$ onto $T \mathbf{V}$ with $\|\operatorname{Proj}\|_{\mathbf{X} \rightarrow T \mathbf{V}}<K$. Now, $T^{-1} \operatorname{Proj} T$ is a projection from $\mathbf{W}$ onto $\mathbf{V}$ of norm less than $D K$.

The type-2 constant of a normed space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ is equal to its "equal norm type-2 constant," namely to the infimum over those $T>0$ for which the second inequality in (136) holds for every $m \in \mathbb{N}$ and every choice of vectors $x_{1}, \ldots, x_{m} \in \mathbf{X}$ that satisfy the additional requirement $\left\|x_{1}\right\|_{\mathbf{X}}=\ldots=\left\|x_{m}\right\|_{\mathbf{X}}$; this is a wellknown result of Pisier, though it first appeared in James' important work [Jam78], where it had a vital role. We will likewise need to use this result, with the twist that we require a small number of unit vectors for which the type- 2 constant of $\mathbf{X}$ is almost attained. The classical proof of the aforementioned equivalence between type-2 and "equal norm type-2" (page 2 of [Jam78]) increases the number of vectors potentially uncontrollably, so we will preform the analysis more carefully in the following lemma, which shows that one need not increase the number of vectors when passing from general vectors to unit vectors.

Lemma 101 (equal norm type 2 without increasing the number of vectors). Fix $n \in \mathbb{N}$ and $0<\beta \leqslant 1$. Let $\left(\mathbf{X},\|\cdot\|_{\mathbf{x}}\right)$ be a normed space and suppose that there exist vectors $x_{1}, \ldots, x_{n} \in \mathbf{X} \backslash\{0\}$ that satisfy

$$
\begin{equation*}
\left(\mathbb{E}\left[\left\|\sum_{i=1}^{n} \mathrm{~g}_{i} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \geqslant \beta T_{2}(\mathbf{X})\left(\sum_{i=1}^{n}\left\|x_{i}\right\|_{\mathbf{X}}^{2}\right)^{\frac{1}{2}} \tag{142}
\end{equation*}
$$

Then, there also exist unit vectors $y_{1}, \ldots, y_{n} \in\left\{x_{i} /\left\|x_{i}\right\|_{\mathbf{X}}\right\}_{i=1}^{n} \subseteq \partial B_{\mathbf{X}}$ that satisfy

$$
\begin{equation*}
\left(\mathbb{E}\left[\left\|\sum_{i=1}^{n} \mathrm{~g}_{i} y_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \gtrsim \beta^{2} T_{2}(\mathbf{X}) \sqrt{n} \tag{143}
\end{equation*}
$$

Proof. We may assume without loss of generality the following normalized version of assumption (142).

$$
\begin{equation*}
\sum_{i=1}^{n}\left\|x_{i}\right\|_{\mathbf{X}}^{2}=1 \quad \text { and } \quad \mathbb{E}\left[\left\|\sum_{i=1}^{n} \mathrm{~g}_{i} x_{i}\right\|_{\mathbf{X}}^{2}\right] \geqslant \beta^{2} T_{2}(\mathbf{X})^{2} \tag{144}
\end{equation*}
$$

For every $k \in \mathbb{N}$ define a subset $I_{k}$ of $\{1, \ldots, n\}$ by

$$
\begin{equation*}
I_{k} \stackrel{\text { def }}{=}\left\{i \in\{1 \ldots, n\}: \frac{1}{2^{k}}<\left\|x_{i}\right\|_{\mathbf{X}} \leqslant \frac{1}{2^{k-1}}\right\} . \tag{145}
\end{equation*}
$$

So, $\left\{I_{k}\right\}_{k \in \mathbb{N}}$ is a partition of $\{1, \ldots, n\}$ as $0<\left\|x_{i}\right\|_{\mathbf{X}} \leqslant 1$ for all $i \in\{1, \ldots, n\}$ by the first equation in (144). Write

$$
\begin{equation*}
m \stackrel{\text { def }}{=}\left[\log _{2}\left(\frac{3 \sqrt{n}}{\beta}\right)\right] \quad \text { and } \quad U \stackrel{\text { def }}{=} \bigcup_{k=1}^{m} I_{k} \times\left\{1, \ldots, 2^{2(m-k)}\right\} \tag{146}
\end{equation*}
$$

With this notation, Lemma 101 will be proven if we show that there exists $S \subseteq U$ with $|S|=n$ such that

$$
\begin{equation*}
\left(\mathbb{E}\left[\left\|\sum_{(i, j) \in S} \frac{\mathrm{~g}_{i j}}{\left\|x_{i}\right\|_{\mathbf{X}}} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \gtrsim \beta^{2} T_{2}(\mathbf{X}) \sqrt{n} \tag{147}
\end{equation*}
$$

where $\left\{\mathrm{g}_{i j}\right\}_{i, j=1}^{\infty}$ are i.i.d. standard Gaussian random variables.
To prove (147), observe first that by the contraction principle (see e.g. [LT91, Section 4.2]) we have

$$
\begin{equation*}
\left(\mathbb{E}\left[\left\|\sum_{(i, j) \in S} \frac{\mathrm{~g}_{i j}}{\left\|x_{i}\right\|_{\mathbf{X}}} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \geqslant\left(\mathbb{E}\left[\left\|\sum_{k=1}^{m} 2^{k-1} \sum_{i \in I_{k}} \sum_{j=1}^{2^{2(m-k)}} \mathbf{1}_{\{(i, j) \in S\}} \mathrm{g}_{i j} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \tag{148}
\end{equation*}
$$

where we used the fact that $1 /\left\|x_{i}\right\|_{\mathbf{X}} \geqslant 2^{k-1}$ for every $k \in \mathbb{N}$ and $i \in I_{k}$ (by the definition (145) of $I_{k}$ ). Also,

$$
1 \stackrel{144}{=} \sum_{i=1}^{n}\left\|x_{i}\right\|_{\mathbf{X}}^{2}=\sum_{k=1}^{\infty} \sum_{i \in I_{k}}\left\|x_{i}\right\|_{\mathbf{X}}^{2} \stackrel{\sqrt[145]{\leqslant}}{\leqslant} \sum_{k=1}^{\infty} \frac{\left|I_{k}\right|}{2^{2 k-2}} \leqslant \frac{4 \sum_{k=1}^{m} 2^{2(m-k)}\left|I_{k}\right|+\sum_{k=m+1}^{\infty}\left|I_{k}\right|}{2^{2 m}} \stackrel{\sqrt[146]{ }}{\leqslant} \frac{\beta^{2}(4|U|+n)}{9 n} .
$$

This simplifies to give that $|U| \geqslant 2 n / \beta^{2}>n$. We can therefore average the right hand side of (148) over all the $n$-point subsets of $U$ a to get the following estimate.

$$
\begin{align*}
\frac{1}{\binom{|U|}{n}} \sum_{\substack{S \subseteq U \\
|S|=n}}\left(\mathbb{E}\left[\left\|\sum_{k=1}^{m} 2^{k-1} \sum_{i \in I_{k}} \sum_{j=1}^{2^{2(m-k)}} \mathbf{1}_{\{(i, j) \in S\}} \mathrm{g}_{i j} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} & \geqslant\left(\mathbb{E}\left[\left\|\sum_{k=1}^{m} 2^{k-1} \sum_{i \in I_{k}} \sum_{j=1}^{2^{2(m-k)}} \frac{\binom{|U|-1}{n-1}}{\binom{|U|}{n}} g_{i j} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \\
& =\frac{n}{2|U|}\left(\mathbb{E}\left[\left\|\sum_{k=1}^{m} 2^{k} \sum_{i \in I_{k}} \sum_{j=1}^{2^{2(m-k)}} \mathrm{g}_{i j} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}}  \tag{149}\\
& =\frac{2^{m-1} n}{|U|}\left(\mathbb{E}\left[\left\|\sum_{k=1}^{m} \sum_{i \in I_{k}} \mathrm{~g}_{i} y_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \\
& \approx \frac{n^{\frac{3}{2}}}{\beta|U|}\left(\mathbb{E}\left[\left\|\sum_{k=1}^{m} \sum_{i \in I_{k}} \mathrm{~g}_{i} y_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}}
\end{align*}
$$

where the first step of (149) uses convexity, the penultimate step of 149 uses the fact that

$$
\left(\left(\sum_{j=1}^{2^{2(m-k)}} \mathrm{g}_{i j}\right)_{i \in I_{k}}\right)_{k=1}^{m} \quad \text { and } \quad\left(\left(2^{m-k} \mathrm{~g}_{i}\right)_{i \in I_{k}}\right)_{k=1}^{m}
$$

have the same distribution, and for the final step of (149) recall the definition (146) of $m$.
It follows from (148) and 149 that there must exist $S \subseteq U$ with $|S|=n$ such that

$$
\begin{equation*}
\left(\mathbb{E}\left[\left\|\sum_{(i, j) \in S} \frac{\mathrm{~g}_{i j}}{\left\|x_{i}\right\|_{\mathbf{X}}} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \gtrsim \frac{n^{\frac{3}{2}}}{\beta|U|}\left(\mathbb{E}\left[\left\|\sum_{k=1}^{m} \sum_{i \in I_{k}} \mathrm{~g}_{i} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \tag{150}
\end{equation*}
$$

To use $(150)$, we claim that $|U| \lesssim n / \beta^{2}$. Indeed,

$$
1 \stackrel{\boxed{144}}{-} \sum_{i=1}^{n}\left\|x_{i}\right\|_{\mathbf{X}}^{2}=\sum_{k=1}^{\infty} \sum_{i \in I_{k}}\left\|x_{i}\right\|_{\mathbf{X}}^{2} \stackrel{\sqrt{145}}{>} \sum_{k=1}^{m} \frac{\left|I_{k}\right|}{2^{2 k}} \stackrel{\mid 146}{-} \frac{|U|}{2^{2 m}} \stackrel{\sqrt[146]{ }}{\geqslant} \frac{\beta^{2}|U|}{81 n} .
$$

By combining the aforementioned upper bound on the size of $U$ with (148) and (150), we see that

$$
\left(\mathbb{E}\left[\left\|\sum_{(i, j) \in S} \frac{\mathrm{~g}_{i j}}{\left\|x_{i}\right\|_{\mathbf{X}}} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \gtrsim \beta \sqrt{n}\left(\mathbb{E}\left[\left\|\sum_{k=1}^{m} \sum_{i \in I_{k}} \mathrm{~g}_{i} y_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}}
$$

From this, we deduce the desired estimate 147 by combining as follows the second inequality in our assumption 144 with the triangle inequality and the definition 136 of the type- 2 constant $T_{2}(\mathbf{X})$.

$$
\begin{aligned}
\left(\mathbb{E}\left[\left\|\sum_{k=1}^{m} \sum_{i \in I_{k}} \mathrm{~g}_{i} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} & \geqslant\left(\mathbb{E}\left[\left\|\sum_{i=1}^{\infty} \mathrm{g}_{i} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}}-\left(\mathbb{E}\left[\left\|\sum_{k=m+1}^{\infty} \sum_{i \in I_{k}} g_{i} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \\
& \stackrel{136}{\geqslant} \beta T_{2}(\mathbf{X})-T_{2}(\mathbf{X})\left(\sum_{k=m+1}^{\infty} \sum_{i \in I_{k}}\left\|x_{i}\right\|_{\mathbf{X}}^{2}\right)^{\frac{1}{2}} \stackrel{145}{\geqslant} \beta T_{2}(\mathbf{X})-\frac{T_{2}(\mathbf{X}) \sqrt{n}}{2^{m}} \stackrel{146}{=} \beta T_{2}(\mathbf{X}) .
\end{aligned}
$$

Proof of Theorem 98 . We will prove that the type 2 constant of $\mathbf{X}$ satisfies

$$
\begin{equation*}
T_{2}(\mathbf{X}) \lesssim \mathrm{LT}(\mathbf{X})(\log n)^{\frac{3}{2}} \tag{151}
\end{equation*}
$$

After (151) will be proven, we deduce Theorem 98 as follows. We first claim that the estimate (151) implies the same upper bound on the cotype 2 constant of $\mathbf{X}$. Namely, we also have

$$
\begin{equation*}
C_{2}(\mathbf{X}) \lesssim \mathrm{LT}(\mathbf{X})(\log n)^{\frac{3}{2}} \tag{152}
\end{equation*}
$$

Indeed,

$$
\begin{equation*}
C_{2}(\mathbf{X}) \leqslant T_{2}\left(\mathbf{X}^{*}\right) \lesssim \mathrm{LT}\left(\mathbf{X}^{*}\right)(\log n)^{\frac{3}{2}} \lesssim \mathrm{LT}(\mathbf{X})(\log n)^{\frac{3}{2}} \tag{153}
\end{equation*}
$$

where the first step of (153) follows from a standard duality argument MP76] (see also e.g. [MS86, Section 9.10], [PW98, Section 4.9] or [AK06, Proposition 6.2.12]), the second step of (153) is an application of (151) to $\mathbf{X}^{*}$, and the third step of (153) is application of Lemma 99 . The desired estimate 139) now follows by a substitution of (151) and (152) into Theorem 91 (Kwapień's theorem).

By [FLM77, Lemma 6.1] (see also the exposition of this fact in [JN10, page 546]) there exists an integer ${ }^{12}$

$$
\begin{equation*}
1 \leqslant m \leqslant \frac{n(n+1)}{2} \tag{154}
\end{equation*}
$$

and $x_{1}, \ldots, x_{m} \in \mathbf{X} \backslash\{0\}$ such that

$$
\begin{equation*}
\left(\mathbb{E}\left[\left\|\sum_{i=1}^{m} \mathrm{~g}_{i} x_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}}=T_{2}(\mathbf{X})\left(\sum_{i=1}^{m}\left\|x_{i}\right\|_{\mathbf{X}}^{2}\right)^{\frac{1}{2}} \tag{155}
\end{equation*}
$$

By Lemma 101, it follows that there exist $y_{1}, \ldots, y_{m} \in \partial B_{\mathbf{X}}$ and a universal constant $0<\gamma<1$ such that

$$
\begin{equation*}
\mathbb{E}\left[\left\|\sum_{i=1}^{m} \mathrm{~g}_{i} y_{i}\right\|_{\mathbf{X}}\right] \geqslant \sqrt{\frac{2}{\pi}}\left(\mathbb{E}\left[\left\|\sum_{i=1}^{m} \mathrm{~g}_{i} y_{i}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}} \geqslant \gamma T_{2}(\mathbf{X}) \sqrt{m} \tag{156}
\end{equation*}
$$

where the first step in (156) holds by (the Gaussian version of) Kahane's inequality Kah64] (see e.g. [LT91, Corollary 3.2] and specifically [LO99, Corollary 3] for the (optimal) constant that we are quoting here even though its value is of secondary importance in the present context). If we denote

$$
\begin{equation*}
\delta \stackrel{\operatorname{def}}{=} \frac{\gamma T_{2}(\mathbf{X})}{\sqrt{m}} \tag{157}
\end{equation*}
$$

then a different way to write 156 is

$$
\begin{equation*}
\mathbb{E}\left[\left\|\sum_{i=1}^{m} \mathrm{~g}_{i} y_{i}\right\|_{\mathbf{X}}\right] \geqslant \delta m \tag{158}
\end{equation*}
$$

Because we ensured that $y_{1}, \ldots, y_{m}$ are unit vectors in $\mathbf{X}$, we may use a theorem of Rudelson and Vershynin [RV06, Theorem 7.4] (an improved Talagrand-style two-parameter version of Elton's theorem; see Remark 102), to deduce from (158) that there are two numbers $0<s \leqslant 1$ and $\delta \lesssim t \leqslant 1$ that satisfy

$$
\begin{equation*}
t \sqrt{s} \gtrsim \frac{\delta}{\left(\log \left(\frac{2}{\delta}\right)\right)^{\frac{3}{2}}} \tag{159}
\end{equation*}
$$

such that there exists a subset $J$ of $\{1, \ldots, m\}$ whose cardinality satisfies

$$
\begin{equation*}
|J| \geqslant s m \tag{160}
\end{equation*}
$$

and moreover we have

$$
\begin{equation*}
\forall\left(a_{j}\right)_{j \in J} \in \mathbb{R}^{J}, \quad t \sum_{j \in J}\left|a_{j}\right| \lesssim\left\|\sum_{j \in J} a_{j} y_{j}\right\|_{\mathbf{X}} \leqslant \sum_{j \in J}\left|a_{j}\right| . \tag{161}
\end{equation*}
$$

(161) means that the Banach-Mazur distance between $\operatorname{span}\left(\left\{y_{j}\right\}_{j \in J}\right)$ and $\ell_{1}^{|J|}$ is $O(1 / t)$. Hence,

$$
\begin{equation*}
\mathrm{c}_{\mathbf{X}}\left(\ell_{1}^{|J|}\right) \lesssim \frac{1}{t} \tag{162}
\end{equation*}
$$

Now, the justification of (151, and hence also the proof of Theorem 98 , can be completed as follows.

$$
\begin{equation*}
\mathrm{LT}(\mathbf{X}) \stackrel{141}{\geqslant} \frac{\mathrm{LT}\left(\ell_{1}^{|J|}\right)}{\mathrm{C}_{\mathbf{X}}\left(\ell_{1}^{|J|}\right)} \stackrel{137)}{\stackrel{162}{\gtrsim}} t \sqrt{|J|} \stackrel{160}{\geqslant} t \sqrt{\operatorname{sm}} \stackrel{159]}{\gtrsim} \frac{\delta \sqrt{m}}{\left(\log \left(\frac{2}{\delta}\right)\right)^{\frac{3}{2}}} \frac{157}{-} \frac{\gamma T_{2}(\mathbf{X})}{\left(\log \left(\frac{2 \sqrt{m}}{\gamma T_{2}(\mathbf{X})}\right)\right)^{\frac{3}{2}}} \gtrsim \frac{T_{2}(\mathbf{X})}{(\log n)^{\frac{3}{2}}} \tag{163}
\end{equation*}
$$

[^8]where the final step of (163) holds because $T_{2}(\mathbf{X}) \geqslant 1$ and $\log m \lesssim \log n$ by (154).
Remark 102. In the proof of Theorem 98 we relied on RV06, Theorem 7.4], which improves (in terms of the power of the logarithm in (159) Talagrand's refinement Tal92] of Elton's theorem Elt83] (which is itself a major quantitative strengthening of an important theorem from (Pis73). Continuing with the notation of Theorem 98, Elton's theorem is a similar statement, except that the size of the subset $J$ is a definite proportion of $m$ that depends only on the parameter $\delta$ for which (158) holds, and also the parameter $t$ for which holds depends only on $\delta$. The asymptotic dependence on $\delta$ in Elton's theorem [Elt83] was improved by Pajor [Paj83], a further improvement was obtained in Tal92], and the optimal dependence on $\delta$ was found by Mendelson and Vershynin in MV03. However, plugging this sharp dependence into our proof of Theorem 98 shows that the classical formulation of Elton's theorem is insufficient for our purposes. The two-parameter formulation of Elton's theorem that was introduced in Ta192] allows for the subset $J$ to have any size through the parameter $s$ in (160), but imposes a relation between $s$ and $t$ such as (159), thus making it possible for us to obtain Theorem 98 .

The only reason why the logarithmic factor in (98) occurs is our use of a Talagrand-style two-parameter version of Elton's theorem, for which the currently best-known bound (RV06) is (159). Thus, if (159) could be improved to $t \sqrt{s} \gtrsim \delta$, i.e., if Question 103 below has a positive answer, then the conclusion (98) of Theorem 98 would become $\mathrm{d}_{\mathrm{x}} \lesssim \mathrm{LT}(\mathbf{X})^{2}$. This would improve Theorem 94 to match the bound of Theorem 93 which is currently known only for infinite dimensional Banach spaces. Moreover, since the resulting bound is independent of the dimension of $\mathbf{X}$, this would yield a new proof of the Lindenstrauss-Tzafriri solution of the complemented subspace problem; the infinite dimensional statement follows formally from its finite dimensional counterpart (e.g. AK06, Theorem 12.1.6]), though all of the steps that led to Theorem 98 work for any reflexive Banach space. Question 103 is interesting in its own right regardless of the above application to the complemented subspace problem. In particular, a positive answer to Question 103 would resolve the question that Talagrand posed in the remark right after Corollary 1.2 in [Tal92], though we warn that he characterises this in [Tal92] as "certainly a rather formidable question."

Question 103. Fix $0<\delta<1$ and $n \in \mathbb{N}$. Let $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ be a Banach space and suppose that $x_{1}, \ldots, x_{n} \in \partial B_{\mathbf{X}}$ satisfy $\mathbb{E}\left[\left\|\sum_{i=1}^{m} \mathrm{~g}_{i} x_{i}\right\|_{\mathbf{X}}\right] \geqslant \delta n$. Does this imply that there are two numbers $0<s, t \leqslant 1$ satisfying $t \sqrt{s} \gtrsim \delta$ and a subset $J \subseteq\{1, \ldots, n\}$ with $|J| \geqslant s n$ such that $\left\|\sum_{j \in J} a_{j} x_{j}\right\|_{\mathbf{X}} \geqslant t \sum_{j \in J}\left|a_{j}\right|$ for every $a_{1}, \ldots, a_{n} \in \mathbb{R}$ ?
2.2. Proof of (106). Because by $\left[\mathrm{CCG}^{+} 98\right]$ we know that $\operatorname{SEP}\left(\ell_{1}^{n}\right)=n$ for every $n \in \mathbb{N}$, using bi-Lipschitz invariance we see that in order to prove (7) it suffices to show that for every normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$,

$$
\begin{equation*}
\exists m \in\{1, \ldots, n\}, \quad \frac{m}{\mathrm{c}_{\mathbf{X}}\left(\ell_{1}^{m}\right)} \geqslant T_{2}(\mathbf{X})^{2} . \tag{164}
\end{equation*}
$$

We will prove (164) using Talagrand's two-parameter refinement of Elton's theorem [Tal92] that we discussed in Remark 102 (the aforementioned improvements over [Tal92] in [MV03, RV06] do not yield a better bound in the ensuing reasoning. Also, the classical formulation of Elton's theorem is insufficient for our purposes, even if one incorporates the asymptotically sharp dependence on $\delta$ from (MV03). Suppose that $k \in \mathbb{N}$ and $x_{1}, \ldots, x_{k} \in B_{\mathbf{X}}$. Let $\mathrm{g}_{1}, \ldots, \mathrm{~g}_{k}$ be i.i.d. standard Gaussian random variables. Denote

$$
E \stackrel{\text { def }}{=} \mathbb{E}\left[\left\|\sum_{j=1}^{k} \mathrm{~g}_{j} x_{j}\right\|_{\mathbf{X}}\right]
$$

By Tal92, Corollary 1.2], there is a universal constant $C \in[1, \infty)$ and a subset $S \subseteq\{1, \ldots, k\}$ satisfying

$$
m \stackrel{\text { def }}{=}|S| \geqslant \frac{E^{2}}{C k},
$$

and such that

$$
\begin{equation*}
\forall\left(a_{j}\right)_{j \in S} \in \mathbb{R}^{S}, \quad \frac{E}{\sqrt{C k m}\left(\log \left(\frac{e C k m}{E^{2}}\right)\right)^{C}} \sum_{j \in S}\left|a_{j}\right| \leqslant\left\|\sum_{j \in S} a_{j} x_{j}\right\|_{\mathbf{X}} \leqslant \sum_{j \in S}\left|a_{j}\right| . \tag{165}
\end{equation*}
$$

Consequently,

$$
\mathrm{c}_{\mathbf{X}}\left(\ell_{1}^{m}\right) \leqslant \frac{\sqrt{C k m}}{E}\left(\log \left(\frac{e C k m}{E^{2}}\right)\right)^{C}
$$

Therefore,

$$
\frac{m}{\mathrm{c}_{\mathbf{X}}\left(\ell_{1}^{m}\right)} \geqslant \frac{E \sqrt{m}}{\sqrt{C k}\left(\log \left(\frac{e C k m}{E^{2}}\right)\right)^{C}} \geqslant \frac{e^{C-\frac{1}{2}}}{2^{C} C^{C+1}} \cdot \frac{E^{2}}{k}=\frac{E^{2}}{k}
$$

where the last step uses the fact that the minimum of the function $u \mapsto \sqrt{u} /\left(\log \left(e C k u / E^{2}\right)\right)^{C}$ on the ray $\left[E^{2} /(C k), \infty\right)$ is attained at $u=e^{2 C-1} E^{2} /(C k)$. It remains to choose $x_{1}, \ldots, x_{k}$ so that $E^{2} / k=T_{2}(\mathbf{X})^{2}$. This is possible because the equal norm type 2 constant of $\mathbf{X}$ equals $T_{2}(\mathbf{X})$, so there are $x_{1}, \ldots, x_{k} \in \partial B_{\mathbf{X}}$ for which

$$
T_{2}(\mathbf{X}) \sqrt{k}=\left(\mathbb{E}\left[\left\|\sum_{j=1}^{k} \mathrm{~g}_{j} x_{j}\right\|_{\mathbf{X}}^{2}\right]\right)^{\frac{1}{2}}=E
$$

where the last step uses Kahane's inequality.
2.3. Hölder extension. In this section we will prove the lower bound on $\mathrm{e}^{\theta}\left(\ell_{\infty}^{n}\right)$ in 20 for every $n \in \mathbb{N}$ and $0<\theta \leqslant 1$. It consists of two estimates, the first of which is

$$
\begin{equation*}
\mathrm{e}^{\theta}\left(\ell_{\infty}^{n}\right) \gtrsim n^{\frac{\theta}{2}+\theta^{2}-1} \tag{166}
\end{equation*}
$$

and the second of which is

$$
\begin{equation*}
\mathrm{e}^{\theta}\left(\ell_{\infty}^{n}\right) \gtrsim n^{\frac{\theta}{4}} . \tag{167}
\end{equation*}
$$

We will justify (166) and (167) separately.
Note that (166) is vacuous if $\theta / 2+\theta^{2}-1 \leqslant 0$, i.e., if $0<\theta \leqslant(\sqrt{17}-1) / 2$. The reason for this is that (166) is based on a reduction to the linear theory from [NR17] (extending the approach of [JL84] to the Hölder regime), that breaks down for functions which are too far from being Lipschitz. Specifically, for a Banach space $\mathbf{X}$ and a closed subspace $\mathbf{E}$ of $\mathbf{X}$, let $\lambda(\mathbf{E} ; \mathbf{X})$ be the projection constant [Grü60] of $\mathbf{E}$ relative to $\mathbf{X}$, i.e., it is the infimum over those $\lambda \in[1, \infty]$ for which there is a projection Proj from $\mathbf{X}$ onto $\mathbf{E}$ whose operator norm satisfies $\|\operatorname{Proj}\|_{\mathbf{X} \rightarrow \mathbf{E}} \leqslant \lambda$. Also, let $\mathrm{e}^{\theta}(\mathbf{X} ; \mathbf{E})$ be the infimum over those $L \in[1, \infty]$ such that for every $\mathcal{C} \subseteq \mathbf{X}$ and every $f: \mathcal{C} \rightarrow \mathbf{E}$ that is $\theta$-Hölder with constant 1 , there is $F: \mathbf{X} \rightarrow \mathbf{E}$ that extend $f$ and is $\theta$-Hölder with constant $L$. With this notation, it was proved in [NR17] (see equation (106) there) that

$$
\begin{equation*}
\mathrm{e}^{\theta}(\mathbf{X} ; \mathbf{E}) \gtrsim \frac{\lambda(\mathbf{E} ; \mathbf{X})^{\theta}}{\operatorname{dim}(\mathbf{E})^{\frac{1-\theta}{2}} \operatorname{dim}(\mathbf{X})^{\theta(1-\theta)} \mathrm{c}_{2}(\mathbf{E})^{1-\theta}} . \tag{168}
\end{equation*}
$$

Using the bounds $\operatorname{dim}(\mathbf{E}) \leqslant \operatorname{dim}(\mathbf{X})$ and $c_{2}(\mathbf{E}) \leqslant \sqrt{\operatorname{dim}(\mathbf{E})}($ John's theorem) in (168), we get that

$$
\begin{equation*}
\mathrm{e}^{\theta}(\mathbf{X} ; \mathbf{E}) \gtrsim \frac{\lambda(\mathbf{E} ; \mathbf{X})^{\theta}}{\operatorname{dim}(\mathbf{X})^{1-\theta^{2}}} . \tag{169}
\end{equation*}
$$

By [Sob41] there is a linear subspace $\mathbf{E}$ of $\ell_{\infty}^{n}$ with $\lambda\left(\mathbf{E} ; \ell_{\infty}^{n}\right)=\sqrt{n}$, using which (169) implies (166).
Remark 104. In NR17] it was deduced from (168) that

$$
\begin{equation*}
\mathrm{e}^{\theta}\left(\ell_{1}^{n}\right) \gtrsim n^{\theta^{2}-\frac{1}{2}} . \tag{170}
\end{equation*}
$$

Specifically, by Kaš77 there is a linear subspace $\mathbf{E}$ of $\ell_{1}^{n}$ with $\mathrm{c}_{2}(\mathbf{E}) \lesssim 1$ and $\operatorname{dim}(\mathbf{E})=\lfloor n / 2\rfloor$; call such $\mathbf{E}$ a Kašin subspace of $\ell_{1}^{n}$. By Rut65] we have $\lambda\left(\mathbf{E} ; \ell_{1}^{n}\right)=\sqrt{n}$, so (170) follows by substituting these parameters into (168). For $\mathbf{X}=\ell_{\infty}^{n}$, the poorly-complemented subspace that we used above can be taken to be the orthogonal complement of any Kašin subspace of $\ell_{1}^{n}$. Such a subspace of $\ell_{\infty}^{n}$ has pathological properties [FJ80]; in particular its Banach-Mazur distance to a Euclidean space is of order $\sqrt{n}$. So, a "vanilla" use of (168) leads at best to (166). However, we expect that it should be possible to improve (166) to

$$
\begin{equation*}
\mathrm{e}^{\theta}\left(\ell_{\infty}^{n}\right) \gtrsim n^{\theta^{2}-\frac{1}{2}} . \tag{171}
\end{equation*}
$$

If (171) holds, then (20) improves to

$$
\mathrm{e}^{\theta}\left(\ell_{\infty}^{n}\right) \gtrsim n^{\max \left\{\frac{\theta}{4}, \theta^{2}-\frac{1}{2}\right\}}=\left\{\begin{array}{lll}
n^{\frac{\theta}{4}} & \text { if } & 0 \leqslant \theta \leqslant \frac{1+\sqrt{33}}{8}  \tag{172}\\
n^{\theta^{2}-\frac{1}{2}} & \text { if } & \frac{1+\sqrt{33}}{8} \leqslant \theta \leqslant 1
\end{array}\right.
$$

For (171), it would suffice to prove the following variant of Conjecture 7 for random subspaces of $\ell_{\infty}^{n}$. Let $\mathbf{E}$ be a subspace of $\mathbb{R}^{n}$ of dimension $m=\lfloor n / 2\rfloor$ that is chosen from the Haar measure on the Grassmannian. We conjecture that there is a universal constant $D \geqslant 1$ such that with high probability there is an origin-symmetric convex body $L \subseteq B_{\mathbf{E}}$ that satisfies $\operatorname{MaxProj}(L) / \operatorname{vol}_{m}(L) \lesssim 1$. If this indeed holds, then by using it in the proof of (168) in [NR17] we can deduce (171) (specifically, replace in Lemma 20 of [NR17] the averaging over $B_{\ell_{2}^{m}}$ by averaging over $L$; we omit the details of this adaptation of [NR17]).

Proof of (167). Fix $k, m \in \mathbb{N}$ satisfying $k \leqslant 2 m \leqslant n / 2$ whose value will be specified later so as to optimize the ensuing reasoning (see (186) below). Denote $\ell=\lfloor(4 m / k)\rfloor$ and define $\mathcal{C}=\mathcal{C}(k, m, n) \subseteq \ell_{\infty}^{n}(\mathbb{C})$ by

$$
\mathcal{C} \stackrel{\text { def }}{=}\left\{E_{m}(k s): s \in\{1, \ldots \ell\}^{n}\right\}
$$

where for every $s=\left(s_{1}, \ldots, s_{n}\right) \in \mathbb{R}^{n}$ we define $E_{m}(s) \in \mathbb{C}^{n}$ by

$$
E_{m}(s) \stackrel{\text { def }}{=} \sum_{j=1}^{n} e^{\frac{\pi i}{2 m} s_{j}} e_{j}
$$

Denote the standard basis (delta masses) of $\mathbb{R}^{\mathcal{C}}$ by $\left\{\boldsymbol{\delta}_{s}\right\}_{s \in \mathcal{C}}$. Let $\mathbb{R}_{0}^{\mathcal{C}}$ be the hyperplane of $\mathbb{R}^{\mathcal{C}}$ consisting of those $\left(a_{s}\right)_{s \in \mathcal{C}}=\sum_{s \in \mathcal{C}} a_{s} \boldsymbol{\delta}_{s}$ with $\sum_{s \in \mathcal{C}} a_{s}=0$. Suppose that $\mathbf{X}_{\theta}=\left(\mathbb{R}_{0}^{\mathcal{C}},\|\cdot\|_{\mathbf{X}_{\theta}}\right)$ is a normed space that satisfies

$$
\begin{equation*}
\forall x, y \in \mathcal{C}, \quad\left\|\boldsymbol{\delta}_{x}-\boldsymbol{\delta}_{y}\right\|_{\mathbf{x}_{\theta}}=\|x-y\|_{\ell_{\infty}^{n}(\mathbb{C})}^{\theta} \tag{173}
\end{equation*}
$$

and,

$$
\begin{equation*}
\forall \mu \in \mathbb{R}_{0}^{\mathcal{C}}, \quad\left(\frac{k}{m}\right)^{\theta}\|\mu\|_{\ell_{1}(\mathcal{C})} \lesssim\|\mu\|_{\mathbf{X}_{\theta}} \lesssim\|\mu\|_{\ell_{1}(\mathcal{C})} \tag{174}
\end{equation*}
$$

For this, $\mathbf{X}_{\theta}$ can be taken to be the normed space whose unit ball is

$$
\begin{equation*}
B_{\mathbf{X}_{\theta}}=\operatorname{conv}\left\{\frac{1}{\|x-y\|_{\ell_{\infty}^{n}(\mathbb{C})}^{\theta}}\left(\boldsymbol{\delta}_{x}-\boldsymbol{\delta}_{y}\right): x, y \in \mathcal{C}, x \neq y\right\} \subseteq \mathbb{R}_{0}^{\mathcal{C}} \tag{175}
\end{equation*}
$$

which is the maximal norm on $\mathbb{R}_{0}^{\complement}$ satisfying (173). To check that (174) holds for the choice (175), note that, as $1 \leqslant k \leqslant 2 m$, distinct $x, y \in \mathcal{C}$ satisfy $k / m \lesssim\|x-y\|_{\ell_{\infty}^{n}(\mathbb{C})} \lesssim 1$. It is simple to deduce (174) from this, as done in [NR17, Lemma 7]. The choice (175) makes $\mathbf{X}_{\theta}$ be the Wasserstein-1 space over ( $\mathcal{C}$, $d_{\theta}$ ), where $d_{\theta}$ is the $\theta$-snowflake of the $\ell_{\infty}^{n}(\mathbb{C})$ metric, i.e., $d_{\theta}(x, y)=\|x-y\|_{\ell_{\infty}^{n}(\mathbb{C})}^{\theta}$ for $x, y \in \ell_{\infty}^{n}(\mathbb{C})$; see Section 5.1 .

By virtue of (173), if we define $f: \mathcal{C} \rightarrow \mathbf{X}_{\theta}$ by setting

$$
\forall x \in \mathcal{C}, \quad f(x) \stackrel{\operatorname{def}}{=} \boldsymbol{\delta}_{x}-\frac{1}{|\mathcal{C}|} \sum_{y \in \mathcal{C}} \boldsymbol{\delta}_{y}
$$

then $f$ is $\theta$-Hölder with constant 1 . We claim that if $m \geqslant \pi \sqrt{n}$, then by 173 every $F: \ell_{\infty}^{n}(\mathbb{C}) \rightarrow \mathbf{X}_{\theta}$ satisfies

$$
\begin{align*}
\frac{1}{(4 m)^{n}} \sum_{j=1}^{n} \sum_{s \in\{1, \ldots, 4 m\}^{n}} & \left\|F\left(E_{m}\left(s+2 m e_{j}\right)\right)-F\left(E_{m}(s)\right)\right\|_{\mathbf{x}_{\theta}} \\
& \lesssim \frac{m^{2+\theta}}{k^{\theta}(12 m)^{n}} \sum_{\varepsilon \in\{-1,0,1\}^{n}} \sum_{s \in\{1, \ldots, 4 m\}^{n}}\left\|F\left(E_{m}(s+\varepsilon)\right)-F\left(E_{m}(s)\right)\right\|_{\mathbf{x}_{\theta}} . \tag{176}
\end{align*}
$$

Indeed, (176) follows from a substitution of (173) into the following inequality from [MN08, Remark 7.5].

$$
\begin{aligned}
& \frac{1}{(4 m)^{n}} \sum_{j=1}^{n} \sum_{s \in\{1, \ldots, 4 m\}^{n}}\left\|F\left(E_{m}\left(s+2 m e_{j}\right)\right)-F\left(E_{m}(s)\right)\right\|_{\ell_{1}(\mathrm{C})} \\
& \lesssim \frac{m^{2}}{(12 m)^{n}} \sum_{\varepsilon \in\{-1,0,1\}^{n}} \sum_{s \in\{1, \ldots, 4 m\}^{n}}\left\|F\left(E_{m}(s+\varepsilon)\right)-F\left(E_{m}(s)\right)\right\|_{\ell_{1}(\mathrm{C})} .
\end{aligned}
$$

Suppose that $F:\{1, \ldots, 4 m\}^{n} \rightarrow \mathbf{X}_{\theta}$ is $\theta$-Hölder with constant $L \geqslant 1$ on $\left(\{1, \ldots, 4 m\}^{n},\|\cdot\|_{\ell_{\infty}^{n}(\mathbb{C})}\right)$, i.e.,

$$
x, y \in\{1, \ldots, 4 m\}^{n}, \quad\|F(x)-F(y)\|_{\mathbf{x}_{\theta}} \leqslant L\|x-y\|_{\ell_{\infty}^{n}(\mathbb{C})}^{\theta}
$$

Then, each of the summands that appear in the right hand side of (176) is at most $2 L / \mathrm{m}^{\theta}$. Consequently,

$$
\begin{equation*}
\frac{1}{n(4 m)^{n}} \sum_{j=1}^{n} \sum_{s \in\{1, \ldots, 4 m\}^{n}}\left\|F\left(E_{m}\left(s+2 m e_{j}\right)\right)-F\left(E_{m}(s)\right)\right\|_{\mathbf{X}_{\theta}} \lesssim \frac{L m^{2}}{k^{\theta} n} \tag{177}
\end{equation*}
$$

If $F$ also extends $f$, then $F\left(E_{m}(s)\right)=f\left(E_{m}\left(s^{\prime}\right)\right)$ for every $s \in \mathbb{N}^{n}$, where $s^{\prime}=\left(s_{1}^{\prime}, \ldots, s_{n}^{\prime}\right)$ and for each $u \in \mathbb{N}$ we let $u^{\prime}$ be an element $\alpha$ of $\{k, 2 k, \ldots, \ell k\}$ for which $|\alpha-u \bmod (4 m)|$ is minimized, so that $s^{\prime} \in \mathcal{C}$ and

$$
\begin{equation*}
\forall s \in \mathbb{N}^{n}, \quad\left\|E_{m}(s)-E_{m}\left(s^{\prime}\right)\right\|_{\ell_{\infty}^{n}(\mathbb{C})} \lesssim \frac{k}{m} \tag{178}
\end{equation*}
$$

Hence, for any $j \in\{1, \ldots, n\}$ and $s \in\{1, \ldots, 4 m\}^{n}$ we have

$$
\begin{align*}
& 2^{\theta}=\left\|-2 e^{\frac{\pi i}{2 m} s_{j}} e_{j}\right\|_{\ell_{\infty}^{n}(\mathbb{C})}^{\theta} \\
&=\left\|E_{m}\left(s+2 m e_{j}\right)-E_{m}(s)\right\|_{\ell_{\infty}^{n}(\mathbb{C})}^{\theta}  \tag{179}\\
& \leqslant\left\|E_{m}\left(\left(s+2 m e_{j}\right)^{\prime}\right)-E_{m}\left(s^{\prime}\right)\right\|_{\ell_{\infty}^{\prime}(\mathbb{C})}^{\theta}+\left\|E_{m}\left(\left(s+2 m e_{j}\right)^{\prime}\right)-E_{m}\left(s+2 m e_{j}\right)\right\|_{\ell_{\infty}^{\prime}(\mathbb{C})}^{\theta}+\left\|E_{m}\left(s^{\prime}\right)-E_{m}(s)\right\|_{\ell_{\infty}^{\prime}(\mathbb{C})}^{\theta} \\
& \leqslant\left\|E_{m}\left(\left(s+2 m e_{j}\right)^{\prime}\right)-E_{m}\left(s^{\prime}\right)\right\|_{\ell_{\infty}^{\prime}(\mathbb{C})}^{\theta}+\frac{2 k^{\theta}}{m^{\theta}}  \tag{180}\\
&=\left\|\boldsymbol{\delta}_{E_{m}\left(\left(s+2 m e_{j}\right)^{\prime}\right)}-\boldsymbol{\delta}_{E_{m}\left(s^{\prime}\right)}\right\|_{\mathbf{X}_{\theta}}+\frac{2 k^{\theta}}{m^{\theta}}  \tag{181}\\
&=\left\|f\left(E_{m}\left(\left(s+2 m e_{j}\right)^{\prime}\right)\right)-f\left(E_{m}\left(s^{\prime}\right)\right)\right\|_{\mathbf{x}_{\theta}}+\frac{2 k^{\theta}}{m^{\theta}}  \tag{182}\\
&=\left\|F\left(E_{m}\left(\left(s+2 m e_{j}\right)^{\prime}\right)\right)-F\left(E_{m}\left(s^{\prime}\right)\right)\right\|_{\mathbf{x}_{\theta}}+\frac{2 k^{\theta}}{m^{\theta}}  \tag{183}\\
& \leqslant\left\|F\left(E_{m}\left(s+2 m e_{j}\right)\right)-F\left(E_{m}(s)\right)\right\|_{\mathbf{X}_{\theta}}+\left\|F\left(E_{m}\left(\left(s+2 m e_{j}\right)^{\prime}\right)\right)-F\left(E_{m}\left(s+2 m e_{j}\right)\right)\right\|_{\mathbf{X}_{\theta}} \\
& \quad+\left\|F\left(E_{m}\left(s^{\prime}\right)\right)-F\left(E_{m}(s)\right)\right\|_{\mathbf{x}_{\theta}}+\frac{2 k^{\theta}}{m^{\theta}} \\
& \leqslant\left\|F\left(E_{m}\left(s+2 m e_{j}\right)\right)-F\left(E_{m}(s)\right)\right\|_{\mathbf{X}_{\theta}}+L\left\|E_{m}\left(\left(s+2 m e_{j}\right)^{\prime}\right)-E_{m}\left(s+2 m e_{j}\right)\right\|_{\ell_{\infty}^{\prime}(\mathbb{C})}^{\theta} \\
& \quad+L\left\|E_{m}\left(s^{\prime}\right)-E_{m}(s)\right\|_{\ell_{\infty}^{n}(\mathbb{C})}^{\theta}+\frac{2 k^{\theta}}{m^{\theta}}  \tag{184}\\
& \leqslant\left\|F\left(E_{m}\left(s+2 m e_{j}\right)\right)-F\left(E_{m}(s)\right)\right\|_{\mathbf{X}_{\theta}}+\frac{2(L+1) k^{\theta}}{m^{\theta}}, \tag{185}
\end{align*}
$$

where for (179) recall the definition of $E_{m}$, in (180) and (185) we used (178), in (181) we used (173), for (182) recall the definition of $f$, in (183) we used the fact that $F$ extends $f$ and $\left\{\left(s+2 m e_{j}\right)^{\prime}, s^{\prime}\right\} \subseteq \mathcal{C}$, and in (184) we used the fact that $F$ is $\theta$-Hölder with constant $L$. By averaging this inequality over ( $j, s$ ) chosen uniformly at random from $\{1, \ldots, n\} \times\{1, \ldots, 4 m\}^{n}$ and applying (177), we conclude that

$$
\begin{equation*}
1 \lesssim\left(\frac{m^{2}}{k^{\theta} n}+\frac{k^{\theta}}{m^{\theta}}\right) L \tag{186}
\end{equation*}
$$

This holds whenever $k, m \in \mathbb{N}$ satisfy $k \leqslant 2 m \leqslant n / 2$ and $m \geqslant \pi \sqrt{n}$, so choose $m=\sqrt{n}$ and $k=\sqrt[4]{n}$ to minimize (up to constants) the right hand side of (186) and deduce the desired lower bound $L \gtrsim n^{\theta / 4}$.

By MN13, Lemma 6.5], for every $\theta \in(0,1]$ and $n \in \mathbb{N}$ we have

$$
\begin{equation*}
\mathrm{e}^{\theta}\left(\ell_{2}^{n}\right) \gtrsim n^{\frac{\theta}{4}} \tag{187}
\end{equation*}
$$

In combination with (167) and AM83], this implies that there is a universal constant $c>0$ such that

$$
\begin{equation*}
\mathrm{e}^{\theta}(\mathbf{X}) \geqslant e^{c \theta \sqrt{\log n}} \tag{188}
\end{equation*}
$$

for every $n$-dimensional normed space $\mathbf{X}$ and every $\theta \in(0,1]$.
Conjecture 105. For any $\theta \in(0,1]$ there is $c(\theta)>0$ such that $\mathrm{e}^{\theta}(\mathbf{X}) \geqslant \operatorname{dim}(\mathbf{X})^{c(\theta)}$ for every normed space $\mathbf{X}$.
Conjecture 105 has a positive answer when the Hölder exponent is close enough to 1 . Specifically, if

$$
\begin{equation*}
0.9307777 \ldots=\frac{\sqrt{193}+1}{16}<\theta \leqslant 1 \tag{189}
\end{equation*}
$$

then

$$
\begin{equation*}
\mathrm{e}^{\theta}(\mathbf{X}) \gtrsim \frac{n^{\frac{\theta\left(8 \theta^{2}-\theta-6\right)}{20 \theta-8}}}{(\log n)^{\frac{3 \theta^{2}}{5 \theta-2}}} \tag{190}
\end{equation*}
$$

Indeed, by bi-Lipschitz invariance, 187 implies the following generalization of Theorem 96 .

$$
\mathrm{e}^{\theta}(\mathbf{X}) \gtrsim \frac{n^{\frac{\theta}{4}}}{\mathrm{~d}_{\mathbf{X}}^{\theta}}
$$

Also,

$$
\mathrm{e}^{\theta}(\mathbf{X}) \stackrel{168}{\gtrsim} \frac{\mathrm{LT}(\mathbf{X})^{\theta}}{n^{(1-\theta)\left(\theta+\frac{1}{2}\right)} \mathrm{d}_{\mathbf{X}}^{1-\theta}} \stackrel{139}{\gtrsim} \frac{\mathrm{~d}_{\mathbf{X}}^{\frac{\theta}{2}} /(\log n)^{\frac{3 \theta}{2}}}{n^{(1-\theta)\left(\theta+\frac{1}{2}\right)} \mathrm{d}_{\mathbf{X}}^{1-\theta}}=\frac{\mathrm{d}_{\mathbf{X}}^{\frac{3 \theta}{2}-1}}{n^{(1-\theta)\left(\theta+\frac{1}{2}\right)}(\log n)^{\frac{3 \theta}{2}}}
$$

Therefore, in analogy to 140 we see that

$$
\begin{equation*}
\mathrm{e}^{\theta}(\mathbf{X}) \gtrsim \max \left\{\frac{n^{\frac{\theta}{4}}}{\mathrm{~d}_{\mathbf{X}}^{\theta}}, \frac{\mathrm{d}_{\mathbf{X}}^{\frac{3 \theta}{2}-1}}{n^{(1-\theta)\left(\theta+\frac{1}{2}\right)}(\log n)^{\frac{3 \theta}{2}}}\right\} \tag{191}
\end{equation*}
$$

Elementary calculus shows that (191) implies (190) in the range (189). If $\theta$ does not satisfy (189), then (191) does not imply a lower bound $\mathrm{e}^{\theta}(\mathbf{X})$ that depends only on $n$ and grows to $\infty$ with $n$; for such $\theta$ the best lower bound that we know is (188). The application of (176) in the above proof of (167) can be mimicked using other bi-Lipschitz invariants to prove (105) for various normed spaces, such as $\ell_{2}^{n}\left(\ell_{1}^{n}\right)$ or $\mathrm{S}_{1}^{n}$, using [NS16] and NS21b], respectively. We do not know if Conjecture 105 holds even when, say, $\mathbf{X}=\ell_{1}^{n}$.
2.4. Justification of (25). In the range $p \in[1,4 / 3] \cup\{2\} \cup[3, \infty]$ the bound in (25) is a combination of [BB12, Corollary 8.12] and [MN13, Theorem 1.17]. We need to justify 25] in the range $p \in(4 / 3,3) \backslash\{2\}$ because it was not previously stated in the literature. Suppose first that $p \in(4 / 3,2)$. By [FLM77], there is $k \in\{1, \ldots, n\}$ with $k=n$ such that $\mathrm{c}_{\ell_{p}^{n}}\left(\ell_{2}^{k}\right) \simeq 1$. Hence,

$$
\mathrm{e}\left(\ell_{p}^{n}\right) \gtrsim \mathrm{e}\left(\ell_{2}^{k}\right) \gtrsim \sqrt[4]{k} \simeq \sqrt[4]{n}
$$

where the penultimate inequality follows from MN13, Theorem 1.17]. Analogously, if $q \in(2,3)$, then by [FLM77] there is $m \in\{1, \ldots, n\}$ with $m=n^{2 / q}$ such that $\mathrm{c}_{\ell_{q}^{n}}\left(\ell_{2}^{m}\right) \simeq 1$. We therefore have

$$
\mathrm{e}\left(\ell_{q}^{n}\right) \gtrsim \mathrm{e}\left(\ell_{2}^{m}\right) \gtrsim \sqrt[4]{m}=n^{\frac{1}{2 q}}
$$

2.5. Proof of the lower bound on $\operatorname{SEP}(\mathbf{X})$ in Theorem 3. Thanks to $\mathbf{7 2}$, the first part of Theorem 106 below coincides with the lower bound on $\operatorname{SEP}(\mathbf{X})$ in Theorem 3 , except that in (192) below we also specify the constant factor that our proof provides (there is no reason to expect that this constant is optimal; due to the fundamental nature of this randomized clustering problem it would be interesting to find the optimal constant here). The second part of Theorem 106 relates to dimension reduction by controlling the cardinality of a finite subset $\mathcal{C}$ of $\mathbf{X}$ on which the lower bound is attained. We conjecture that the first part of 193 below could be improved to $|\mathcal{C}|^{1 / n}=O(1)$; an inspection of the ensuing proof suggests that a possible route towards this improved bound is to incorporate a proportional Dvoretzky-Rogers factorization [BS88, ST89, Gia96] in place of our use of the "vanilla" Dvoretzky-Rogers lemma [DR50].
Theorem 106. For every $n \in \mathbb{N}$, any n-dimensional normed space $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right)$ satisfies

$$
\begin{equation*}
\operatorname{SEP}(\mathbf{X}) \geqslant \operatorname{evr}(\mathbf{X}) \frac{2(n!)^{\frac{1}{2 n}} \Gamma\left(1+\frac{n}{2}\right)^{\frac{1}{n}}}{\sqrt{\pi n}}=\frac{\sqrt{2}+o(1)}{e \sqrt{\pi}} \operatorname{evr}(\mathbf{X}) \sqrt{n} . \tag{192}
\end{equation*}
$$

Furthermore, there exists a finite subset $\bigodot$ of $\mathbf{X}$ satisfying

$$
\begin{equation*}
|\mathcal{C}|^{\frac{1}{n}} \lesssim \frac{\sqrt{n}}{\operatorname{evr}(\mathbf{X})} \quad \text { and } \quad \operatorname{SEP}\left(\mathcal{C}_{\mathbf{X}}\right) \gtrsim \operatorname{evr}(\mathbf{X}) \sqrt{n} \tag{193}
\end{equation*}
$$

Our proof of Theorem 106 builds upon the strategy that was used in $\left[\mathrm{CG}^{+} 98\right]$ to treat $\ell_{1}^{n}$. A combinatorial fact on which it relies is Lemma 107 below, which is implicit in the proof of [CCG ${ }^{+} 98$, Lemma 3.1]. After proving Theorem 106 while using Lemma 107 , we will present a proof of Lemma 107 which is a quick application of the Loomis-Whitney inequality [LW49]; the proof in [CCG ${ }^{+} 98$ uses a result of AKPW91] which is proved in AKPW91 via information-theoretic reasoning through the use of Shearer's inequality [CGFS86]; the relation between the Loomis-Whitney inequality and Shearer's inequality is well-known (see e.g. $\overline{\mathrm{BB} 12]}$ ), so our proof of Lemma 107 is in essence a repackaging of the classical ideas.
Lemma 107. Fix $n, M \in \mathbb{N}$ and a nonempty finite subset $\Omega$ of $\mathbb{Z}^{n}$. Suppose that $\mathcal{P}$ is a random partition of $\Omega$ that is supported on partitions into subsets of cardinality at most $M$, i.e.,

$$
\operatorname{Prob}\left[\max _{\Gamma \in \mathcal{P}}|\Gamma| \leqslant M\right]=1
$$

Then, there exists $i \in\{1, \ldots, n\}$ and $x \in \Omega \cap\left(\Omega-e_{i}\right)$ for which

$$
\begin{equation*}
\operatorname{Prob}\left[\mathcal{P}(x) \neq \mathcal{P}\left(x+e_{i}\right)\right] \geqslant \frac{1}{\sqrt[n]{M}}-\frac{1}{n} \sum_{i=1}^{n} \frac{\left|\Omega \backslash\left(\Omega-e_{i}\right)\right|}{|\Omega|} \tag{194}
\end{equation*}
$$

Proof of Theorem 106 assuming Lemma 107 . By suitably choosing the identification of $\mathbf{X}$ with $\mathbb{R}^{n}$, we may assume without loss of generality that $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ and $B_{\ell_{2}^{n}}$ is the Löwner ellipsoid of $B_{\mathbf{X}}$. Then,

$$
\begin{equation*}
\operatorname{evr}(\mathbf{X})=\left(\frac{\operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}\right)^{\frac{1}{n}}=\frac{\sqrt{\pi}}{\Gamma\left(1+\frac{n}{2}\right)^{\frac{1}{n}} \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{\frac{1}{n}}} \tag{195}
\end{equation*}
$$

By the Dvoretzky-Rogers lemma [DR50], there exist contact points $x_{1}, \ldots, x_{n} \in S^{n-1} \cap \partial B_{\mathbf{X}}$ that satisfy

$$
\begin{equation*}
\forall k \in\{1, \ldots, n\}, \quad\left\|\operatorname{Proj}_{\operatorname{span}\left(x_{1}, \ldots, x_{k-1}\right)^{\perp}}\left(x_{k}\right)\right\|_{\ell_{2}^{n}} \geqslant \sqrt{\frac{n-k+1}{n}} . \tag{196}
\end{equation*}
$$

Let $\Lambda=\Lambda\left(x_{1}, \ldots, x_{n}\right) \subseteq \mathbb{R}^{n}$ denote the lattice that is generated by $x_{1}, \ldots, x_{n}$, namely

$$
\Lambda=\sum_{i=1}^{n} \mathbb{Z} x_{i}=\left\{\sum_{i=1}^{n} k_{i} x_{i}: k_{1}, \ldots, k_{n} \in \mathbb{Z}\right\}
$$

By (196), $\Lambda$ is full-rank. Denote the fundamental parallelepiped of $\Lambda$ by $Q=Q\left(x_{1}, \ldots, x_{n}\right)$, i.e.,

$$
Q=\sum_{i=1}^{n}[0,1) x_{i}=\left\{\sum_{i=1}^{n} s_{i} x_{i}: 0 \leqslant s_{1}, \ldots, s_{n}<1\right\} .
$$

Since $x_{1}, \ldots, x_{n} \in B_{\mathbf{X}}$, we have $Q-Q \subseteq n B_{\mathbf{X}}$ and by (196) the volume of $Q$ (the determinant of $\Lambda$ ) satisfies

$$
\begin{equation*}
\operatorname{det}(\Lambda)=\operatorname{vol}_{n}(Q)=\prod_{k=1}^{n}\left\|\operatorname{Proj}_{\operatorname{span}\left\{x_{1}, \ldots, x_{k-1}\right\}^{\perp}}\left(x_{k}\right)\right\|_{\ell_{2}^{n}} \stackrel{\sqrt{196}}{\geqslant} \prod_{k=1}^{n} \sqrt{\frac{n-k+1}{n}}=\frac{\sqrt{n!}}{n^{\frac{n}{2}}} . \tag{197}
\end{equation*}
$$

Fix $m \in \mathbb{N}$ and $\sigma, \Delta>0$. Denote $\mathcal{C}_{m}=\mathcal{C}_{m}\left(x_{1}, \ldots, x_{n}\right)=\Lambda \cap(m Q)=\left\{\sum_{i=1}^{n} k_{i} x_{i}: k_{1}, \ldots, k_{n} \in\{0, \ldots, m-1\}\right\}$ and suppose that $\mathcal{P}$ is $\sigma$-separating $\Delta$-bounded random partition of $\mathcal{C}_{m}$. The $\Delta$-boundedness of $\mathcal{P}$ means that $\Gamma-\Gamma \subseteq \Delta B_{\mathbf{X}}$ for every $\Gamma \subseteq \mathcal{C}_{m}$ with $\operatorname{Prob}[\Gamma \in \mathcal{P}]>0$. Recalling that $Q-Q \subseteq n B_{\mathbf{X}}$, this implies that

$$
\begin{equation*}
B \mathbf{X} \supseteq \frac{1}{\Delta+n}((\Gamma+Q)-(\Gamma+Q)) . \tag{198}
\end{equation*}
$$

Now,

$$
\begin{equation*}
\frac{\sqrt{\pi}}{\Gamma\left(1+\frac{n}{2}\right)^{\frac{1}{n}} \operatorname{evr}(\mathbf{X})}=\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{\frac{1}{n}} \geqslant \frac{2}{\Delta+n} \operatorname{vol}_{n}(\Gamma+Q)^{\frac{1}{n}}=\frac{2}{\Delta+n}\left(|\Gamma| \operatorname{vol}_{n}(Q)\right)^{\frac{1}{n}} \geqslant \frac{2(n!)^{\frac{1}{2 n}}}{(\Delta+n) \sqrt{n}}|\Gamma|^{\frac{1}{n}}, \tag{199}
\end{equation*}
$$

where the first step of (199) is (195), the second step of (199) uses (198) and the Brunn-Minkowski inequality, the third step of (199) holds because the parallelepipeds $\{\gamma+Q: \gamma \in \Gamma\}$ are disjoint, and the final step of (199) is (197). If $T \in \mathrm{GL}_{n}(\mathbb{R})$ is given by $T e_{i}=x_{i}$, then it follows from (199) that the random partition $T^{-1} \mathcal{P}=\left\{T^{-1} \Gamma: \Gamma \in \mathcal{P}\right\}$ of $T^{-1} \mathcal{C}_{m}=\{0, \ldots, m-1\}^{n}$ satisfies the assumptions of Lemma 107 with

$$
M=\frac{(\pi n)^{\frac{n}{2}}(\Delta+n)^{n}}{2^{n} \Gamma\left(1+\frac{n}{2}\right) \sqrt{n!}} \cdot \frac{1}{\operatorname{evr}(\mathbf{X})^{n}} .
$$

If we choose $\Omega=\{0, \ldots, m-1\}^{n}=T^{-1} \mathcal{C}_{m}$ in Lemma 107 then $|\Omega|=m^{n}$ and $\left|\Omega \backslash\left(\Omega-e_{i}\right)\right|=m^{n-1}$ for every $i \in\{1, \ldots, n\}$, so it follows from Lemma 107 that there exist $i \in\{1, \ldots, n\}$ and $x \in \mathcal{C}_{m}$ such that

$$
\begin{equation*}
\operatorname{Prob}\left[\mathcal{P}(x) \neq \mathcal{P}\left(x+x_{i}\right)\right] \geqslant \operatorname{evr}(\mathbf{X}) \frac{2(n!)^{\frac{1}{2 n}} \Gamma\left(1+\frac{n}{2}\right)^{\frac{1}{n}}}{(\Delta+n) \sqrt{\pi n}}-\frac{1}{m} . \tag{200}
\end{equation*}
$$

At the same time, the left hand side of (200) is at most $\sigma / \Delta$, since $\mathcal{P}$ is $\sigma$-separating and $\left\|x_{i}\right\|_{\mathbf{X}} \leqslant 1$. Thus,

$$
\begin{equation*}
\sigma \geqslant \operatorname{evr}(\mathbf{X}) \frac{2 \Delta(n!)^{\frac{1}{2 n}} \Gamma\left(1+\frac{n}{2}\right)^{\frac{1}{n}}}{(\Delta+n) \sqrt{\pi n}}-\frac{\Delta}{m} . \tag{201}
\end{equation*}
$$

By letting $m \rightarrow \infty$ in (201) and then letting $\Delta \rightarrow \infty$ in the resulting estimate, we get (192). Also, if we set $\Delta=n$ in 201, then for sufficiently large $m=\sqrt{n} / \operatorname{evr}(\mathbf{X})$ we have $\operatorname{SEP}\left(\mathcal{C}_{m}\right) \gtrsim \operatorname{evr}(\mathbf{X}) \sqrt{n}$, giving (193).

We will next provide a proof of Lemma 107 whose main ingredient is the following lemma.
Lemma 108 (application of Loomis-Whitney). Fix an integer $n \geqslant 2$ and a finite subset $\Gamma$ of $\mathbb{Z}^{n}$. For $x \in \mathbb{Z}^{n}$ and $i \in\{1, \ldots, n\}$, let $d_{i}(x ; \Gamma) \in \mathbb{N} \cup\{0\}$ be the number of times that the oriented discrete axis-parallel line $x+\mathbb{Z} e_{i}$ transitions from $\Gamma$ to $\mathbb{Z}^{n} \backslash \Gamma$, and let $\mathrm{g}(x ; \Gamma)$ be the geometric mean of $d_{1}(x ; \Gamma), \ldots, d_{n}(x ; \Gamma)$. Thus

$$
\forall i \in\{1, \ldots, n\}, \quad d_{i}(x ; \Gamma) \stackrel{\text { def }}{=}\left|\left\{k \in \mathbb{Z}: x+k e_{i} \in \Gamma \wedge x+(k+1) e_{i} \notin \Gamma\right\}\right|,
$$

and

$$
g(x ; \Gamma) \stackrel{\text { def }}{=} \sqrt[n]{d_{1}(x ; \Gamma) \cdots d_{n}(x ; \Gamma)}
$$

Then,

$$
\begin{equation*}
\frac{1}{n} \sum_{i=1}^{n}\left|\Gamma \backslash\left(\Gamma-e_{i}\right)\right| \geqslant\left(\sum_{x \in \mathbb{Z}^{n}} g(x ; \Gamma)^{\frac{n}{n-1}}\right)^{\frac{n-1}{n}} \geqslant|\Gamma|^{\frac{n-1}{n}} . \tag{202}
\end{equation*}
$$

Proof. The second inequality in (202) holds because $d_{1}(x ; \Gamma), \ldots, d_{n}(x ; \Gamma) \geqslant 1$ for every $x \in \Gamma$ (as $\left.|\Gamma|<\infty\right)$, and hence $g(\cdot ; \Gamma) \geqslant \mathbf{1}_{\Gamma}(\cdot)$ point-wise. For the first inequality in [202), observe that for each $i \in\{1, \ldots, n\}$,

$$
\left|\Gamma \backslash\left(\Gamma-e_{i}\right)\right|=\sum_{x \in \mathbb{Z}^{n}} \mathbf{1}_{\Gamma}(x) \mathbf{1}_{\mathbb{Z}^{n} \backslash \Gamma}\left(x+e_{i}\right)=\sum_{y \in \operatorname{Proj}_{e_{i}} \Gamma}\left(\sum_{k \in \mathbb{Z}} \mathbf{1}_{\Gamma}\left(y+k e_{i}\right) \mathbf{1}_{\mathbb{Z}^{n} \backslash \Gamma}\left(y+(k+1) e_{i}\right)\right)=\sum_{y \in \operatorname{Proj}_{e_{i}} \mathbb{Z}^{n}} d_{i}(y ; \Gamma) .
$$

Consequently,
$\frac{1}{n} \sum_{i=1}^{n}\left|\Gamma \backslash\left(\Gamma-e_{i}\right)\right|=\frac{1}{n} \sum_{i=1}^{n}\left\|d_{i}(\cdot ; \Gamma)^{\frac{1}{n-1}}\right\|_{\ell_{n-1}\left(\operatorname{Proj}_{e_{i}} \mathbb{Z}^{n}\right)}^{n-1} \geqslant \prod_{i=1}^{n}\left\|d_{i}(\cdot ; \Gamma)^{\frac{1}{n-1}}\right\|_{\ell_{n-1}\left(\operatorname{Proj}_{e_{i}} \mathbb{Z}^{n}\right)}^{\frac{n-1}{n}} \geqslant \sum_{x \in \mathbb{Z}^{n}} \prod_{i=1}^{n} d_{i}\left(\operatorname{Proj}_{e_{i}^{i}} x\right)^{\frac{1}{n-1}}$,
where the second step is an application of the arithmetic-mean/geometric-mean inequality and the final step is an application of the Loomis-Whitney inequality [LW49] (see [Sil73], Theorem 3] for the functional version of the Loomis-Whitney inequality that the are using here); we note that even though this inequality is commonly stated for functions on $\mathbb{R}^{n}$ rather than for functions on $\mathbb{Z}^{n}$, its proof for functions on $\mathbb{Z}^{n}$ is identical (in fact, [LW49] proves the continuous inequality by first proving its discrete counterpart).

Note that when $n=1$ Lemma 108 holds trivially if we interpret (202) as $|\Gamma \backslash(\Gamma-1)| \geqslant \max _{x \in \mathbb{Z}} g(x ; \Gamma) \geqslant 1$, since in this case $g(x ; \Gamma)=|\Gamma \backslash(\overline{\Gamma-1)})|$ for every $x \in \mathbb{Z}$.

The following corollary of Lemma 108 is a deterministic counterpart of Lemma 107 .
Corollary 109. Fix $n, M \in \mathbb{N}$ and a nonempty finite subset $\Omega$ of $\mathbb{Z}^{n}$. Suppose that $\mathcal{P}$ is a partition of $\Omega$ with

$$
\begin{equation*}
\max _{\Gamma \in \mathcal{P}}|\Gamma| \leqslant M . \tag{203}
\end{equation*}
$$

Then,

$$
\begin{equation*}
\frac{1}{n} \sum_{i=1}^{n}\left|\left\{x \in \Omega \cap\left(\Omega-e_{i}\right): \mathcal{P}(x) \neq \mathcal{P}\left(x+e_{i}\right)\right\}\right| \geqslant \frac{|\Omega|}{\sqrt[n]{M}}-\frac{1}{n} \sum_{i=1}^{n}\left|\Omega \backslash\left(\Omega-e_{i}\right)\right| \tag{204}
\end{equation*}
$$

Proof. Observe that for each fixed $i \in\{1, \ldots, n\}$ we have

$$
\begin{align*}
\left|\Omega \backslash\left(\Omega-e_{i}\right)\right|+\sum_{x \in \Omega \cap\left(\Omega-e_{i}\right)} \mathbf{1}_{\mathcal{P}(x) \neq \mathcal{P}\left(x+e_{j}\right)} & =\left|\Omega \backslash\left(\Omega-e_{i}\right)\right|+\sum_{x \in \Omega \cap\left(\Omega-e_{i}\right)}\left(\sum_{\Gamma \in \mathcal{P}} \mathbf{1}_{\Gamma}(x) \mathbf{1}_{\mathbb{Z}^{n} \backslash \Gamma}\left(x+e_{i}\right)\right) \\
& =\sum_{x \in \mathbb{Z}^{n}} \sum_{\Gamma \in \mathcal{P}} \mathbf{1}_{\Gamma}(x) \mathbf{1}_{\mathbb{Z}^{n} \backslash \Gamma}\left(x+e_{i}\right)  \tag{205}\\
& =\sum_{\Gamma \in \mathcal{P}}\left|\Gamma \backslash\left(\Gamma-e_{i}\right)\right|,
\end{align*}
$$

where the first step of 205 holds because $\mathcal{P}$ is a partition of $\Omega$ and the second step of (205) holds because $\mathbf{1}_{\Gamma}(x) \mathbf{1}_{\mathbb{Z}^{n} \backslash \Gamma}\left(x+e_{i}\right)=0$ for every $\Gamma \subseteq \Omega$ if $x \in \mathbb{Z}^{n} \backslash \Omega$, and if $x \in \Omega \backslash\left(\Omega-e_{i}\right)$, then $\mathbf{1}_{\Gamma}(x) \mathbf{1}_{\mathbb{Z}^{n}} \backslash \Gamma\left(x+e_{i}\right)=1$ for exactly one $\Gamma \in \mathcal{P}$ (specifically, this holds for $\Gamma=\mathcal{P}(x)$ because $x+e_{i} \in \mathbb{Z}^{n} \backslash \Omega \subseteq \mathbb{Z}^{n} \backslash \mathcal{P}(x)$ ). Now,

$$
\begin{aligned}
& \frac{1}{n} \sum_{i=1}^{n}\left|\left\{x \in \Omega \cap\left(\Omega-e_{i}\right): \mathcal{P}(x) \neq \mathcal{P}\left(x+e_{i}\right)\right\}\right|+\frac{1}{n} \sum_{i=1}^{n}\left|\Omega \backslash\left(\Omega-e_{i}\right)\right| \\
& \quad \stackrel{[205]}{=} \sum_{\Gamma \in \mathcal{P}} \frac{1}{n} \sum_{i=1}^{n}\left|\Gamma \backslash\left(\Gamma-e_{i}\right)\right| \stackrel{[202]}{\geqslant} \sum_{\Gamma \in \mathcal{P}}|\Gamma|^{\frac{n-1}{n}} \stackrel{[203]}{\geqslant} \frac{1}{\sqrt[n]{M}} \sum_{\Gamma \in \mathcal{P}}|\Gamma|=\frac{|\Omega|}{\sqrt[n]{M}},
\end{aligned}
$$

where the last step holds because $\mathcal{P}$ is a partition of $\Omega$.
Proof of Lemma 107. Denoting $p=\max _{i \in\{1, \ldots, n\}} \max _{x \in \Omega \cap\left(\Omega-e_{i}\right)} \operatorname{Prob}\left[\mathcal{P}(x) \neq \mathcal{P}\left(x+e_{i}\right)\right]$, the goal is to show that $p$ is at least the right hand side of (194). This follows from Corollary 109 because

$$
\begin{aligned}
p|\Omega| \geqslant \frac{p}{n} \sum_{i=1}^{n}\left|\Omega \cap\left(\Omega-e_{i}\right)\right| \geqslant \frac{1}{n} \sum_{i=1}^{n} \sum_{x \in \Omega \cap\left(\Omega-e_{i}\right)} \operatorname{Prob}\left[\mathcal{P}(x) \neq \mathcal{P}\left(x+e_{i}\right)\right]=\frac{1}{n} \sum_{i=1}^{n} \sum_{x \in \Omega \cap\left(\Omega-e_{i}\right)} \mathbb{E}\left[\mathbf{1}_{\mathcal{P}(x) \neq \mathcal{P}\left(x+e_{i}\right)}\right] \\
=\mathbb{E}\left[\frac{1}{n} \sum_{i=1}^{n}\left|\left\{x \in \Omega \cap\left(\Omega-e_{i}\right): \mathcal{P}(x) \neq \mathcal{P}\left(x+e_{i}\right)\right\}\right|\right] \stackrel{\mid 204}{\geqslant} \frac{|\Omega|}{\sqrt[n]{M}}-\frac{1}{n} \sum_{i=1}^{n}\left|\Omega \backslash\left(\Omega-e_{i}\right)\right| .
\end{aligned}
$$

2.6. Proof of the lower bound on $\operatorname{PAD}_{\delta}(\mathbf{X})$ in Theorem 68 . Fixing $n \in \mathbb{N}$, a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$, and $\delta \in(0,1)$, recalling the notation in Definition 64 we will prove here that

$$
\begin{equation*}
\operatorname{PAD}_{\delta}(\mathbf{X}) \geqslant \sup _{m \in \mathbb{N}} \mathrm{PAD}_{\delta}^{m}(\mathbf{X}) \geqslant \frac{2}{1-\sqrt[n]{\delta}} \tag{206}
\end{equation*}
$$

which gives the first inequality in 103.

Proof of 206). Suppose that $0<\varepsilon<1$ and $r>2$. Let $\eta_{\varepsilon}$ be any $\varepsilon$-net of $r B_{\mathbf{X}}$. Then, $\log \left|\eta_{\varepsilon}\right|=n \log (r / \varepsilon)$ (see e.g. Ost13, Lemma 9.18]). Fix a (disjoint) Voronoi tessellation $\left\{V_{x}\right\}_{x \in n_{\varepsilon}}$ of $r B_{\mathbf{X}}$ that is induced by $n_{\varepsilon}$. Thus, $\left\{V_{x}\right\}_{x \in \eta_{\varepsilon}}$ is a partition of $r B_{\mathbf{X}}$ into Borel subsets such that $x \in V_{x} \subseteq x+\varepsilon B_{\mathbf{X}}$ for every $x \in n_{\varepsilon}$. So, for every $w \in r B_{\mathbf{X}}$ there is a unique net point $x(w) \in n_{\varepsilon}$ such that $w \in V_{x(w)}$.

Fix $\mathfrak{p}>\sup _{m \in \mathbb{N}} \operatorname{PAD}_{\delta}^{m}(\mathbf{X}) \geqslant \operatorname{PAD}_{\delta}\left(n_{\varepsilon}\right)$, and assume from now on that $0<\varepsilon<1 /(2 \mathfrak{p})$ and $r>1 / \mathfrak{p}-2 \varepsilon$ (eventually we will consider the limits $\varepsilon \rightarrow 0$ and $r \rightarrow \infty$ ). By the definition of $\mathrm{PAD}_{\delta}\left(n_{\varepsilon}\right)$, there exists a probability distribution $\mathcal{P}$ over 1-bounded partitions of $\eta_{\varepsilon}$ such that

$$
\begin{equation*}
\forall y \in n_{\varepsilon}, \quad \operatorname{Prob}\left[\left(y+\frac{1}{\mathfrak{p}} B_{\mathbf{X}}\right) \cap n_{\varepsilon} \subseteq \mathcal{P}(y)\right] \geqslant \delta \tag{207}
\end{equation*}
$$

For every $y \in n_{\varepsilon}$ define

$$
\mathcal{P}^{*}(y) \stackrel{\text { def }}{=} \bigcup_{z \in \mathcal{P}(y)} V_{z}=\left\{w \in r B_{\mathbf{X}}: x(w) \in \mathcal{P}(y)\right\}
$$

Then $\left\{\mathcal{P}^{*}(y)\right\}_{y \in n_{\varepsilon}}$ is a (finitely supported) random partition of $r B_{\mathbf{X}}$ into Borel subsets.
We claim that for every $y \in n_{\varepsilon}$ the following inclusion of events holds.

$$
\begin{equation*}
\left\{w \in \mathbb{R}^{n}: w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(y)\right\}+\frac{1-2 \varepsilon \mathfrak{p}}{(1+2 \varepsilon) \mathfrak{p}}\left(\mathcal{P}^{*}(y)-\mathcal{P}^{*}(y)\right) \subseteq \mathcal{P}^{*}(y) \tag{208}
\end{equation*}
$$

Indeed, take any $w \in \mathbb{R}^{n}$ such that

$$
w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(y)
$$

and also take any $u, v \in \mathcal{P}^{*}(y)$. By the definition of $\mathcal{P}^{*}$ we have $x(u), x(\nu) \in \mathcal{P}(y)$. As $\mathcal{P}$ is 1-bounded, we have $\|x(u)-x(v)\|_{\mathbf{X}} \leqslant 1$. Therefore $\|u-v\|_{\mathbf{X}} \leqslant\|u-x(u)\|_{\mathbf{X}}+\|x(u)-x(v)\|_{\mathbf{X}}+\|v-x(v)\|_{\mathbf{X}} \leqslant 1+2 \varepsilon$. Hence,

$$
\frac{1-2 \varepsilon \mathfrak{p}}{(1+2 \varepsilon) \mathfrak{p}}(u-v) \in \frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}}
$$

so the assumption on $w$ implies that

$$
w+\frac{1-2 \varepsilon \mathfrak{p}}{(1+2 \varepsilon) \mathfrak{p}}(u-v) \in \mathcal{P}^{*}(y)
$$

This is precisely the assertion in (208). By the Brunn-Minkowski inequality, 208) gives

$$
\operatorname{vol}_{n}\left(\mathcal{P}^{*}(y)\right)^{\frac{1}{n}} \geqslant 2 \frac{1-2 \varepsilon \mathfrak{p}}{(1+2 \varepsilon) \mathfrak{p}} \operatorname{vol}_{n}\left(\mathcal{P}^{*}(y)\right)^{\frac{1}{n}}+\operatorname{vol}_{n}\left(\left\{w \in \mathbb{R}^{n}: w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(y)\right\}\right)^{\frac{1}{n}}
$$

This simplifies to give the following estimate.

$$
\begin{equation*}
\operatorname{vol}_{n}\left(\left\{w \in \mathbb{R}^{n}: w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(y)\right\}\right) \leqslant\left(1-2 \frac{1-2 \varepsilon \mathfrak{p}}{(1+2 \varepsilon) \mathfrak{p}}\right)^{n} \operatorname{vol}_{n}\left(\mathcal{P}^{*}(y)\right) \tag{209}
\end{equation*}
$$

Now,

$$
\begin{align*}
& \operatorname{vol}_{n}\left(\left\{w \in r B_{\mathbf{X}}: w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(x(w))\right\}\right) \\
&=\sum_{y \in n_{\varepsilon}} \operatorname{vol}_{n}\left(\left\{w \in \mathcal{P}^{*}(y): w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(x(w))\right\}\right)  \tag{210}\\
&=\sum_{y \in n_{\varepsilon}} \operatorname{vol}_{n}\left(\left\{w \in \mathcal{P}^{*}(y): w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(y)\right\}\right)  \tag{211}\\
& \leqslant\left(1-2 \frac{1-2 \varepsilon \mathfrak{p}}{(1+2 \varepsilon) \mathfrak{p}}\right)^{n} \sum_{y \in n_{\varepsilon}} \operatorname{vol}_{n}\left(\mathcal{P}^{*}(y)\right)  \tag{212}\\
&=\left(1-2 \frac{1-2 \varepsilon \mathfrak{p}}{(1+2 \varepsilon) \mathfrak{p}}\right)^{n} r^{n} \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) . \tag{213}
\end{align*}
$$

Here (210) holds because $\left\{\mathcal{P}^{*}(y)\right\}_{y \in n_{\varepsilon}}$ is a partition of $r B_{\mathbf{X}}$. The identity (211) holds because, since by the definition of $\mathcal{P}^{*}$ we have $w \in \mathcal{P}^{*}(x(w))$ for every $w \in r B_{\mathbf{X}}$ and the sets $\left\{\mathcal{P}^{*}(y)\right\}_{y \in n_{\varepsilon}}$ are pairwise disjoint, if $w \in \mathcal{P}^{*}(y)$ for some $y \in n_{\varepsilon}$ then necessarily $\mathcal{P}^{*}(x(w))=\mathcal{P}^{*}(y)$. The estimate 212 uses 209. The identity (213) uses once more that $\left\{\mathcal{P}^{*}(y)\right\}_{y \in n_{\varepsilon}}$ is a partition of $r B_{\mathbf{X}}$.

We next claim that for every $w \in(r+2 \varepsilon-1 / \mathfrak{p}) B_{\mathbf{X}}$ the following inclusion of events holds.

$$
\begin{equation*}
\left\{\left(x(w)+\frac{1}{\mathfrak{p}} B_{\mathbf{X}}\right) \cap n_{\varepsilon} \subseteq \mathcal{P}(x(w))\right\} \subseteq\left\{w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(x(w))\right\} . \tag{214}
\end{equation*}
$$

Indeed, suppose that $w \in \mathbf{X}$ satisfies $\|w\|_{\mathbf{X}} \leqslant r+2 \varepsilon-1 / \mathfrak{p}$ and $\left(x(w)+(1 / \mathfrak{p}) B_{\mathbf{X}}\right) \cap n_{\varepsilon} \subseteq \mathcal{P}(x(w))$. Fix any $z \in \mathbf{X}$ such that $\|w-z\|_{\mathbf{X}} \leqslant(1-2 \varepsilon \mathfrak{p}) / \mathfrak{p}$. Then we have $\|z\|_{\mathbf{X}} \leqslant\|w\|_{\mathbf{X}}+\|w-z\|_{\mathbf{X}} \leqslant r$, so $z \in r B_{\mathbf{X}}$ and therefore $x(z) \in n_{\varepsilon}$ is well-defined. Now,

$$
\|x(w)-x(z)\|_{\mathbf{X}} \leqslant\|x(w)-w\|_{\mathbf{X}}+\|w-z\|_{\mathbf{X}}+\|z-x(z)\|_{\mathbf{X}} \leqslant \varepsilon+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}}+\varepsilon=\frac{1}{\mathfrak{p}}
$$

Hence, our assumption on $w$ implies that $x(z) \in \mathcal{P}(x(w))$. By the definition of $\mathcal{P}^{*}(x(w))$, this means that $z \in \mathcal{P}^{*}(x(w)$ ), thus completing the verification of (214). Due to (207) and (214) we conclude that

$$
\begin{equation*}
\forall w \in\left(r+2 \varepsilon-\frac{1}{\mathfrak{p}}\right) B_{\mathbf{X}}, \quad \operatorname{Prob}\left[w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(x(w))\right] \geqslant \delta . \tag{215}
\end{equation*}
$$

Finally,

$$
\begin{aligned}
\delta\left(r+2 \varepsilon-\frac{1}{\mathfrak{p}}\right)^{n} \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) & \stackrel{215}{\leqslant} \int_{\left(r+2 \varepsilon-\frac{1}{\mathfrak{p}}\right) B_{\mathbf{X}}} \operatorname{Prob}\left[w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(x(w))\right] \mathrm{d} w \\
& =\mathbb{E}\left[\operatorname{vol}_{n}\left(\left\{w \in\left(r+2 \varepsilon-\frac{1}{\mathfrak{p}}\right) B_{\mathbf{X}}: w+\frac{1-2 \varepsilon \mathfrak{p}}{\mathfrak{p}} B_{\mathbf{X}} \subseteq \mathcal{P}^{*}(x(w))\right\}\right)\right] \\
& \stackrel{213}{\leqslant}\left(1-2 \frac{1-2 \varepsilon \mathfrak{p}}{(1+2 \varepsilon) \mathfrak{p}}\right)^{n} r^{n} \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) .
\end{aligned}
$$

This simplifies to give the estimate

$$
\sqrt[n]{\delta}\left(1-\frac{1}{\mathfrak{p} r}+\frac{2 \varepsilon}{r}\right) \leqslant 1-2 \frac{1-2 \varepsilon \mathfrak{p}}{(1+2 \varepsilon) \mathfrak{p}} .
$$

By letting $r \rightarrow \infty$, then $\varepsilon \rightarrow 0$, and then $\mathfrak{p} \rightarrow \sup _{m \in \mathbb{N}} \operatorname{PAD}_{\delta}^{m}(\mathbf{X})$, the desired bound (206) follows.
2.7. Proof of Proposition 86. The final lower bound from the Introduction that remains to be proven is Proposition 86. The ensuing reasoning is a restructuring of a proof that was shown to us by Lutwak.

Lemma 110. Every origin-symmetric convex body $K \subseteq \mathbb{R}^{n}$ satisfies

$$
\begin{equation*}
\int_{S^{n-1}} \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{u^{\perp}}(K)\right)}{\|u\|_{K}^{n+1}} \mathrm{~d} u \geqslant \frac{n^{2} \Gamma\left(\frac{n}{2}\right)}{2 \sqrt{\pi} \Gamma\left(\frac{n+1}{2}\right)} \operatorname{vol}_{n}(K)^{2} . \tag{216}
\end{equation*}
$$

Equality in 216) holds if and only if $K$ is an ellipsoid.
Before proving Lemma 110 , we will explain how it implies Proposition 86
Proof of Proposition 86assuming Lemma 110. The following standard identity follows from integration in polar coordinates (its quick derivation can be found, for example, on page 91 of [Pis89]).

$$
\begin{equation*}
\operatorname{vol}_{n}(K)=\frac{1}{n} \int_{S^{n-1}} \frac{\mathrm{~d} u}{\|u\|_{K}^{n}} . \tag{217}
\end{equation*}
$$

Hence,

$$
\begin{align*}
& \int_{S^{n-1}} \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{u^{\perp}}(K)\right)}{\|u\|_{K}^{n+1}} \mathrm{~d} u \leqslant\left(\int_{S^{n-1}} \frac{\mathrm{~d} u}{\|u\|_{K}^{n}}\right) \max _{u \in S^{n-1}} \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{u^{\perp}}(K)\right)}{\|u\|_{K}}  \tag{218}\\
& \quad \stackrel{\boxed{217}}{=} n \operatorname{vol}_{n}(K) \max _{z \in \partial K}\left(\|z\|_{2}^{n} \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z^{\perp}}(K)\right)\right)=n^{2} \operatorname{vol}_{n}(K) \max _{z \in \partial K} \operatorname{vol}_{n}\left(\operatorname{Cone}_{z}(K)\right) .
\end{align*}
$$

The desired inequality (129) follows by contrasting (218) with (216). Consequently, if there is equality in (129), then (216) must hold as equality as well, so the characterization of the equality case in Proposition 86 follows from the characterization of the quality case in Lemma 110 .

The important Petty projection inequality [Pet71] (see also [Sch95, MM96] for different proofs, as well as the survey (Lut93]) states that for every convex body $K \subseteq \mathbb{R}^{n}$, the affine invariant quantity

$$
\begin{equation*}
\operatorname{vol}_{n}(K)^{n-1} \operatorname{vol}_{n}\left(\Pi^{*} K\right) \tag{219}
\end{equation*}
$$

is maximized when $K$ is an ellipsoid, and ellipsoids are the only maximizers of (219). Recall that the polar projection body $\Pi^{*} K$ is given by (30), which shows in particular that $\operatorname{vol}_{n-1}\left(B_{\ell_{2}^{n-1}}\right) \Pi^{*} B_{\ell_{2}^{n}}=B_{\ell_{2}^{n}}$. Hence,

$$
\operatorname{vol}_{n}(K)^{n-1} \operatorname{vol}_{n}\left(\Pi^{*} K\right) \leqslant \operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right)^{n-1} \operatorname{vol}_{n}\left(\Pi^{*} B_{\ell_{2}^{n}}\right)=\left(\frac{\operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right)}{\operatorname{vol}_{n-1}\left(B_{\ell_{2}^{n-1}}\right)}\right)^{n}=\left(\frac{2 \sqrt{\pi} \Gamma\left(\frac{n+1}{2}\right)}{n \Gamma\left(\frac{n}{2}\right)}\right)^{n} .
$$

At the same time, by combining (30) and (217) we have

$$
\operatorname{vol}_{n}\left(\Pi^{*} K\right)=\frac{1}{n} \int_{S^{n-1}} \frac{\mathrm{~d} u}{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{u^{\perp}}(K)\right)^{n}} .
$$

Consequently, Petty's projection inequality can be restated as the following estimate,

$$
\begin{equation*}
\int_{S^{n-1}} \frac{\mathrm{~d} u}{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{u^{\perp}}(K)\right)^{n}} \leqslant\left(\frac{2 \sqrt{\pi} \Gamma\left(\frac{n+1}{2}\right)}{n \Gamma\left(\frac{n}{2}\right)}\right)^{n} \frac{n}{\operatorname{vol}_{n}(K)^{n-1}} \tag{220}
\end{equation*}
$$

together with the assertion that holds as an equality if and only if $K$ is an ellipsoid.

Proof of Lemma 110. Observe that

$$
\begin{align*}
\operatorname{vol}_{n}(K) & =\frac{1}{n} \int_{S^{n-1}}\left(\frac{1}{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{u^{\perp}}(K)\right)^{\frac{n}{n+1}}}\right)\left(\frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{u^{\perp}}(K)\right)^{\frac{n}{n+1}}}{\|u\|_{K}^{n}}\right) \mathrm{d} u  \tag{221}\\
& \leqslant \frac{1}{n}\left(\int_{S^{n-1}} \frac{\mathrm{~d} u}{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{u^{\perp}}(K)\right)^{n}}\right)^{\frac{1}{n+1}}\left(\int_{S^{n-1}} \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{u^{\perp}}(K)\right)}{\|u\|_{K}^{n+1}} \mathrm{~d} u\right)^{\frac{n}{n+1}}  \tag{222}\\
& \leqslant \frac{1}{n}\left(\frac{2 \sqrt{\pi} \Gamma\left(\frac{n+1}{2}\right)}{n \Gamma\left(\frac{n}{2}\right)}\right)^{\frac{n}{n+1}} \frac{n^{\frac{1}{n+1}}}{\operatorname{vol}_{n}(K)^{\frac{n-1}{n+1}}}\left(\int_{S^{n-1}} \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{u^{\perp}}(K)\right)}{\|u\|_{K}^{n+1}} \mathrm{~d} u\right)^{\frac{n}{n+1}}, \tag{223}
\end{align*}
$$

where (221) is (217), in (222) we used Hölder's inequality with the conjugate exponents $1+\frac{1}{n}$ and $n+1$, and (223) is an application of (220). This simplifies to give the desired inequality (216).
Remark 111. Fix $n \in \mathbb{N}$, a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ and $x \in S^{n-1}$. Both of the bounds in (50) follow from elementary geometric reasoning (convexity and Fubini's theorem). Recalling (30), the second inequality in (50) is $\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}} B_{\mathbf{X}}\right) \leqslant n\|x\|_{\mathbf{X}} \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) / 2$; its justification can be found in the proof of Lemma 5.1 in [GNS12] (this was not included in the version of [GNS12] that appeared in the journal, but it appears in the arxiv version of [GNS12]). The rest of (50) is vol $_{n}\left(B_{\mathbf{X}}\right)\|x\|_{\mathbf{X}} \leqslant 2 \mathrm{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}} B_{\mathbf{X}}\right)$; since we did not find a reference for the derivation of this simple lower bound on hyperplane projections, we will now quickly justify it. For every $u \in \operatorname{Proj}_{x^{\perp}} B_{\mathbf{X}}$ let $s(u)=\inf \left\{s \in \mathbb{R}: u+s x \in B_{\mathbf{X}}\right\}$ and $t(u)=\sup \left\{t \in \mathbb{R}: u+t x \in B_{\mathbf{X}}\right\}$. For every $u \in \operatorname{Proj}_{x^{\perp}} B_{\mathbf{X}}$ we have $u+t(u) x \in B_{\mathbf{X}}$, and by symmetry also $-u-s(u) x \in B_{\mathbf{X}}$. Hence, by convexity

$$
\frac{1}{2}(u+t(u) x)+\frac{1}{2}(-u-s(u) x)=\frac{t(u)-s(u)}{2} x \in B_{\mathbf{X}} .
$$

By the definition of $t(0)$, this means that $(t(u)-s(u)) / 2 \leqslant t(0)=1 /\|x\|_{\mathbf{X}}$. Consequently, using Fubini's theorem (recall that $x \in S^{n-1}$ ) we conclude that

$$
\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)=\int_{\operatorname{Pro}_{x^{\prime}} \perp B_{\mathbf{X}}}(t(u)-s(u)) \mathrm{d} u \leqslant \int_{\operatorname{Proj}_{x} \perp B_{\mathbf{X}}} \frac{2}{\|x\|_{\mathbf{X}}} \mathrm{d} u=\frac{2}{\|x\|_{\mathbf{X}}} \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}} B_{\mathbf{X}}\right) .
$$

## 3. Preliminaries on random partitions

This section treats basic properties of random partitions, including measurability issues that we need for subsequent applications. As such, it is of a technical/foundational nature and it can be skipped on first reading if one is willing to accept the measurability requirements that are used in the proofs that appear in Section 4 and Section 5 .

Recall that a random partition $\mathcal{P}$ of a metric space $(m, d m$ ) was defined in the Introduction as follows. One is given a probability space ( $\Omega, \mathbf{P r o b}$ ) and a sequence of set-valued mappings $\left\{\Gamma^{k}: \Omega \rightarrow 2^{m}\right\}_{k=1}^{\infty}$ such that for each fixed $k \in \mathbb{N}$ the mapping $\Gamma^{k}: \Omega \rightarrow 2^{m}$ is strongly measurable relative to the $\sigma$-algebra of Prob-measurable subsets of $\Omega$, i.e., the set $\left(\Gamma^{k}\right)^{-}(E)=\left\{\omega \in \Omega: E \cap \Gamma^{k}(\omega) \neq \varnothing\right\}$ is Prob-measurable for every closed $E \subseteq M$. We require that $\mathcal{P}^{\omega}=\left\{\Gamma^{k}(\omega)\right\}_{k=1}^{\infty}$ is a partition of $M$ for every $\omega \in \Omega$.

Definition 62 and Definition 64 (of separating and padded random partitions, respectively) assumed implicitly that the quantities that appear in the left hand sides of equations (93) and (96) are well-defined, i.e., that the events $\{\mathcal{P}(x) \neq \mathcal{P}(y)\}$ and $\{B m(x, r) \subseteq \mathcal{P}(x)\}$ are Prob-measurable for every $x, y \in m$ and $r>0$. This follows from the above definition, because for every closed subset $E \subseteq M$ we have

$$
\left\{\omega \in \Omega: \mathcal{P}^{\omega}(x) \neq \mathcal{P}^{\omega}(y)\right\}=\bigcup_{\substack{k, \ell \in \in \mathbb{N} \\ k \neq \ell}}\left(\left\{\omega \in \Omega:\{x\} \cap \Gamma^{k}(\omega) \neq \varnothing\right\} \cap\left\{\omega \in \Omega:\{y\} \cap \Gamma^{\ell}(\omega) \neq \varnothing\right\}\right),
$$

and

$$
\left\{\omega \in \Omega: E \nsubseteq \mathcal{P}^{\omega}(x)\right\}=\bigcup_{\substack{k, \ell \in \mathbb{N} \\ k \neq \ell}}\left(\left\{\omega \in \Omega:\{x\} \cap \Gamma^{k}(\omega) \neq \varnothing\right\} \cap\left\{\omega \in \Omega: E \cap \Gamma^{\ell}(\omega) \neq \varnothing\right\}\right) .
$$

Another "leftover" from the Introduction is the proof of Lemma 66, which asserts that the moduli of Definition 62 and Definition 64 are bi-Lipschitz invariants. The proof of this simple but needed statement is the following direct use of the definition of a $\Delta$-bounded random partition.

Proof of Lemma 66. Fix $D>\mathrm{c}_{\left(n, d_{n}\right)}\left(m, d_{m}\right)$. There is an embedding $\phi: m \rightarrow n$ and a scaling factor $\lambda>0$ such that (16) holds. Fix $\Delta>0$ and let $\mathcal{P}$ be a $\lambda \Delta$-bounded random partition of $n$. Suppose that $\mathcal{P}$ is induced by the probability space ( $\Omega$, Prob), i.e., there are strongly measurable mappings $\left\{\Gamma^{k}: \Omega \rightarrow 2^{n}\right\}_{k=1}^{\infty}$ such that $\mathcal{P}^{\omega}=\left\{\Gamma^{k}(\omega)\right\}_{k=1}^{\infty}$ for every $\omega \in \Omega$. For every $k \in \mathbb{N}$ the mapping $\omega \mapsto \phi^{-1}\left(\Gamma^{k}(\omega)\right) \in 2^{m}$ is strongly measurable. Indeed, if $E \subseteq W$ is closed then, because $M$ is complete and $\phi$ is a homeomorphism, also $\phi(E) \subseteq \eta$ is closed. So, $\left\{\omega \in \Omega: \phi(E) \cap \Gamma^{k}(\omega) \neq \varnothing\right\}=\left\{\omega \in \Omega: E \cap \phi^{-1}\left(\Gamma^{k}(\omega)\right) \neq \varnothing\right\}$ is Prob-measurable, as required. Therefore, if we define $Q^{\omega}=\left\{\phi^{-1}\left(\Gamma^{k}(\omega)\right)\right\}_{k=1}^{\infty}$ for $\omega \in \Omega$, then $Q$ is a random partition of $m$.
$Q$ is $\Delta$-bounded because for every $x \in T$ and $u, v \in \mathcal{Q}(x)$ we have $\phi(u), \phi(v) \in \mathcal{P}(\phi(x))$, and therefore $d_{m}(u, v) \leqslant d_{n}(\phi(u), \phi(\nu)) / \lambda \leqslant \operatorname{diam}_{\eta}(\mathcal{P}(\phi(x))) / \lambda \leqslant \Delta$, using (16) and that $\mathcal{P}$ is $\lambda \Delta$-bounded. For every $x, y \in W$ the events $\{\mathcal{Q}(x) \neq \mathscr{Q}(y)\}$ and $\{\mathcal{P}(\phi(x)) \neq \mathcal{P}(\phi(y))\}$ coincide. So, if $\mathcal{P}$ is $\sigma$-separating for some $\sigma>0$,

$$
\operatorname{Prob}[\mathcal{L}(x) \neq \mathcal{Q}(y)]=\operatorname{Prob}[\mathcal{P}(\phi(x)) \neq \mathcal{P}(\phi(y))] \leqslant \frac{\sigma}{\lambda \Delta} d_{n}(\phi(x), \phi(y)) \stackrel{\sqrt[16]{ }}{\stackrel{D \sigma}{\leqslant}} \frac{D \sigma}{\Delta} d_{m}(x, y)
$$

This shows that Q is $(D \sigma)$-separating, thus establishing the first assertion (98) of Lemma66.
Suppose that $\mathcal{P}$ is $(\mathfrak{p}, \delta)$-padded for some $\mathfrak{p}>0$ and $0<\delta<1$. Fix $x \in M$. Assuming that the event $\left\{B_{n}(\phi(x), \lambda \Delta / \mathfrak{p}) \subseteq \mathcal{P}(\phi(x))\right\}$ occurs, if $z \in B_{m}(x, \Delta /(D \mathfrak{p}))$, then $d_{\eta}(\phi(z), \phi(x)) \leqslant \lambda D d_{m}(z, x) \leqslant \lambda \Delta / \mathfrak{p}$ by (16). Thus, $\phi(z) \in B_{\eta}(\phi(x), \lambda \Delta / \mathfrak{p})$ and therefore $\phi(z) \in \mathcal{P}(\phi(x))$, i.e., $z \in \mathcal{Q}(x)$. This shows the inclusion of events $\left\{B_{n}(\phi(x), \lambda \Delta / \mathfrak{p}) \subseteq \mathcal{P}(\phi(x))\right\} \subseteq\left\{B_{m}(x, \Delta /(D \mathfrak{p})) \subseteq \mathcal{Q}(x)\right\}$. Since $\mathcal{P}$ is $(\mathfrak{p}, \delta)$-padded, it follows from this that also $Q$ is $(D \mathfrak{p}, \delta)$-padded, thus establishing the second assertion (99) of Lemma66.

The final basic "leftover" from the Introduction is the following simple proof of Lemma67.
Proof of Lemma 67. Fix $\Delta>0$ and suppose that $\sigma_{1}>\operatorname{SEP}\left(m_{1}\right)$ and $\sigma_{2}>\operatorname{SEP}\left(m_{2}\right)$. Define

$$
\begin{equation*}
\Delta_{1}=\Delta\left(\frac{\sigma_{1}}{\sigma_{1}+\sigma_{2}}\right)^{\frac{1}{s}} \quad \text { and } \quad \Delta_{2}=\Delta\left(\frac{\sigma_{2}}{\sigma_{1}+\sigma_{2}}\right)^{\frac{1}{s}} . \tag{224}
\end{equation*}
$$

Let $\mathcal{P}_{\Delta_{1}}$ be a $\sigma_{1}$-separating $\Delta_{1}$-bounded random partition of $m_{1}$. Similarly, let $\mathcal{P}_{\Delta_{2}}$ be a $\sigma_{2}$-separating $\Delta_{2}$-bounded random partition of $m_{2}$. Assume that $\mathcal{P}_{\Delta_{1}}$ and $\mathcal{P}_{\Delta_{2}}$ are independent random variables. Let $\mathcal{P}_{\Delta}$ be the corresponding product random partition of $m_{1} \times m_{2}$, i.e., its clusters are give by

$$
\begin{equation*}
\forall\left(x_{1}, x_{2}\right) \in m_{1} \times m_{2}, \quad \mathcal{P}_{\Delta}\left(x_{1}, x_{2}\right)=\mathcal{P}_{\Delta_{1}}\left(x_{1}\right) \times \mathcal{P}_{\Delta_{2}}\left(x_{2}\right) . \tag{225}
\end{equation*}
$$

By (224) we have $\Delta_{1}^{s}+\Delta_{2}^{s}=\Delta^{s}$, so $\mathcal{P}_{\Delta}$ is a $\Delta$-bounded random partition of $m_{1} \oplus_{s} m_{2}$ (the required measurability is immediate). It remains to note that every $\left(x_{1}, x_{2}\right),\left(y_{1}, y_{2}\right) \in m_{1} \times m_{2}$ satisfy

$$
\begin{align*}
\operatorname{Prob}\left[\mathcal{P}_{\Delta}\left(x_{1}, x_{2}\right) \neq \mathcal{P}_{\Delta}\left(y_{1}, y_{2}\right)\right] & =1-\operatorname{Prob}\left[\mathcal{P}_{\Delta_{1}}\left(x_{1}\right)=\mathcal{P}_{\Delta_{1}}\left(y_{1}\right)\right] \operatorname{Prob}\left[\mathcal{P}_{\Delta_{2}}\left(x_{2}\right)=\mathcal{P}_{\Delta_{2}}\left(y_{2}\right)\right]  \tag{226}\\
& \leqslant 1-\left(1-\frac{\sigma_{1} d_{m_{1}}\left(x_{1}, y_{1}\right)}{\Delta_{1}}\right)\left(1-\frac{\sigma_{2} d_{m_{2}}\left(x_{2}, y_{2}\right)}{\Delta_{2}}\right)  \tag{227}\\
& =\frac{\sigma_{1} d m_{1}\left(x_{1}, y_{1}\right)}{\Delta_{1}}+\frac{\sigma_{2} d_{m_{2}}\left(x_{2}, y_{2}\right)}{\Delta_{2}}-\frac{\sigma_{1} \sigma_{2} d m_{1}\left(x_{1}, y_{1}\right) d_{m_{2}}\left(x_{2}, y_{2}\right)}{\Delta_{1} \Delta_{2}}  \tag{228}\\
& \leqslant\left(\left(\frac{\sigma_{1}}{\Delta_{1}}\right)^{\frac{s}{s-1}}+\left(\frac{\sigma_{2}}{\Delta_{2}}\right)^{\frac{s}{s-1}}\right)^{\frac{s-1}{s}}\left(d_{m_{1}}\left(x_{1}, y_{1}\right)^{s}+d_{m_{2}}\left(x_{2}, y_{2}\right)^{s}\right)^{\frac{1}{s}}  \tag{229}\\
& =\frac{\sigma_{1}+\sigma_{2}}{\Delta} d_{m_{1} \oplus} m_{2}\left(\left(x_{1}, x_{2}\right),\left(y_{1}, y_{2}\right)\right), \tag{230}
\end{align*}
$$

where (226) uses (225) and the independence of $\mathcal{P}_{\Delta_{1}}$ and $\mathcal{P}_{\Delta_{2}}$, the bound (227) is an application of the assumption that $\mathcal{P}_{\Delta_{1}}$ is $\sigma_{1}$-separating and $\mathcal{P}_{\Delta_{2}}$ is $\sigma_{2}$-separating, (229) is an application of Hölder's inequality, and (230) follows from (100) and (224). This proves (101). Note that even though we dropped
the quadratic additive improvement in (228), this does not change the final bound in (101) due to the need to work with all possible scales $\Delta>0$ and all possible values of $d_{m_{1}}\left(x_{1}, y_{1}\right)$ and $d_{m_{2}}\left(x_{2}, y_{2}\right)$.

To prove (102), fix $\mathfrak{p}_{1}>\operatorname{PAD}_{\delta_{1}}\left(m_{1}\right)$ and $\mathfrak{p}_{2}>\operatorname{PAD}_{\delta_{2}}\left(m_{2}\right)$ and replace (224) by

$$
\Delta_{1}=\frac{\Delta \mathfrak{p}_{1}}{\left(\mathfrak{p}_{1}^{s}+\mathfrak{p}_{2}^{s}\right)^{\frac{1}{s}}} \quad \text { and } \quad \Delta_{2}=\frac{\Delta \mathfrak{p}_{2}}{\left(\mathfrak{p}_{1}^{s}+\mathfrak{p}_{2}^{s}\right)^{\frac{1}{s}}} .
$$

This time, we choose $\mathcal{P}_{\Delta_{1}}$ to be a $\left(\mathfrak{p}_{1}, \delta_{1}\right)$-padded $\Delta_{1}$-bounded random partition of $m_{1}$. Similarly, let $\mathcal{P}_{\Delta_{2}}$ be a $\left(\mathfrak{p}_{2}, \delta_{2}\right)$-padded $\Delta_{2}$-bounded random partition of $m_{2}$, with $\mathcal{P}_{\Delta_{1}}$ and $\mathcal{P}_{\Delta_{2}}$ independent, and we again combine them as in 225) to give the product partition $\mathcal{P}_{\Delta}$ of $m_{1} \times m_{2}$. The analogous reasoning shows that $\mathcal{P}_{\Delta}$ is a $\left(\left(\mathfrak{p}_{1}^{s}+\mathfrak{p}_{2}^{s}\right)^{1 / s}, \delta_{1} \delta_{2}\right)$-padded $\Delta$-bounded random partition of $m_{1} \oplus_{s} m_{2}$.
3.1. Standard set-valued mappings. Recall that a metric space $\left(m, d_{m}\right)$ is said to be Polish if it is separable and complete. Polish metric spaces are the appropriate setting for Lipschitz extension theorems that are based on the assumption that for every $\Delta>0$ there is a probability distribution over $\Delta$-bounded partitions of $m$ with certain properties. Indeed, a Banach space-valued Lipschitz function can always be extended to the completion of $m$ while preserving the Lipschitz constant, and the mere existence of countably many sets of diameter at most $\Delta$ that cover $m$ for every $\Delta>0$ implies that $m$ is separable.

Theorem 65 assumes local compactness. Even though this assumption is more restrictive than being Polish, it suffices for the applications that we obtain herein because they deal with finite dimensional normed spaces. It is, however, possible to treat general Polish metric spaces by working with a notion of measurability of set-valued mappings that differs from the strong measurability that was assumed in Section 1.7. We call this notion standard set-valued mappings; see Definition 112 .

The requirements for a set-valued mapping to be standard are quite innocuous and easy to check. In particular, the clusters of the specific random partitions that we will study are easily seen to be standard set-valued mappings. It is also simple to verify that the clusters of the random partitions that we construct are strongly measurable. So, we have two approaches, which are both easy to work with. We chose to work in the Introduction with the requirement that the clusters are strongly measurable because this directly makes the quantity $\operatorname{SEP}(\cdot)$ be bi-Lipschitz invariant, and it is also slightly simpler to describe. Nevertheless, in practice it is straightforward to check that the clusters are standard, and even though we do not know that this leads to a bi-Lipschitz invariant (we suspect that it does not), it does lead to an easily implementable Lipschitz extension criterion that holds in the maximal generality of Polish spaces.

Definition 112 (standard set-valued mapping). Suppose that $\left(\mathscr{L}, d_{\mathcal{L}}\right)$ is a Polish metric space and that $\Omega \subseteq \mathscr{L}$ is a Borel subset of $\mathscr{L}$. Given a metric space ( $m, d_{m}$ ), a set-valued mapping $\Gamma: \Omega \rightarrow 2^{m}$ is said to be standard if the following three conditions hold.

- For every $x \in T$ the set $\{\omega \in \Omega: x \in \Gamma(\omega)\}$ is Borel.
- The set $\mathcal{G}_{\Gamma}=\Gamma^{-}(m)=\{\omega \in \Omega: \Gamma(\omega) \neq \varnothing\}$ is Borel.
- For every $x \in T$ the mapping $\left(\omega \in \mathcal{G}_{\Gamma}\right) \mapsto d_{m}(x, \Gamma(\omega))$ is Borel measurable on $\mathcal{G}_{\Gamma}$.

The following extension criterion is a counterpart to Theorem 65 that works in the maximal generality of Polish metric spaces; its proof, which is an adaptation of ideas of [LN05], appears in Section 5

Theorem 113. Let $\left(m, d_{m}\right)$ be a Polish metric space and fix another metric $\mathfrak{d}$ on $M$. Suppose that for every $\Delta>0$ there is a Polish metric space $\mathscr{L}_{\Delta}$, a Borel subset $\Omega_{\Delta} \subseteq \mathscr{L}_{\Delta}$, a Borel probability measure Prob $_{\Delta}$ on $\Omega_{\Delta}$ and a sequence of standard set-valued mappings $\left\{\Gamma_{\Delta}^{k}: \Omega_{\Delta} \rightarrow 2^{m}\right\}_{k=1}^{\infty}$ such that $\mathcal{P}_{\Delta}^{\omega}=\left\{\Gamma_{\Delta}^{k}(\omega)\right\}_{k=1}^{\infty}$ is a partition of $M$ for every $\omega \in \Omega_{\Delta}$, for every $x \in M$ and $\omega \in \Omega_{\Delta}$ we have $\operatorname{diam} m\left(\mathcal{P}_{\Delta}^{\omega}(x)\right) \leqslant \Delta$, and

$$
\begin{equation*}
\forall x, y \in M, \quad \Delta \operatorname{Prob}_{\Delta}\left[\omega \in \Omega_{\Delta}: \mathcal{P}_{\Delta}^{\omega}(x) \neq \mathcal{P}_{\Delta}^{\omega}(y)\right] \leqslant \mathfrak{d}(x, y) \tag{231}
\end{equation*}
$$

Then, for every Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$, every subset $\mathcal{C} \subseteq M$ and every 1 -Lipschitz mapping $f: \mathcal{C} \rightarrow \mathbf{Z}$, there exists a mapping $F: M \rightarrow \mathbf{Z}$ that extends $f$ and satisfies $\|F(x)-F(y)\|_{\mathbf{z}} \lesssim \mathfrak{d}(x, y)$ for every $x, y \in M$ (namely, $F$ is Lipschitz on $T$ with respect to the metric $\mathfrak{d}$ ). Moreover, $F$ depends linearly on $f$.
3.2. Proximal selectors. For later applications we need to know that set-valued mappings that are either strongly measurable or standard admit certain auxiliary measurable mappings that are (perhaps approximately) the closest point to a given (but arbitrary) nonempty closed subset of the metric space in question. We will justify this now using classical descriptive set theory.

Lemma 114. Fix a measurable space $(\Omega, \mathcal{F})$. Suppose that $(M, d m)$ is a metric space and that $S \subseteq M$ is nonempty and locally compact. Let $\Gamma: \Omega \rightarrow 2^{m}$ be a strongly measurable set-valued mapping such that $\Gamma(\omega)$ is a bounded subset of $m$ for every $\omega \in \Omega$. Then there exists an $\mathcal{F}$-to-Borel measurable mapping $\gamma: \Omega \rightarrow S$ that satisfies $d m(\gamma(\omega), \Gamma(\omega))=d_{m}(S, \Gamma(\omega))$ for every $\omega \in \Omega$ for which $\Gamma(\omega) \neq \varnothing$.

Proof. For every $\omega \in \Omega$ define a subset $\Phi(\omega) \subseteq S$ as follows.

$$
\Phi(\omega) \stackrel{\text { def }}{=}\left\{\begin{array}{cl}
\left\{s \in S: d_{m}(s, \Gamma(\omega))=d_{m}(S, \Gamma(\omega))\right\} & \text { if } \Gamma(\omega) \neq \varnothing \\
S & \text { if } \Gamma(\omega)=\varnothing
\end{array}\right.
$$

The goal of Lemma 114 is to demonstrate the existence of an $\mathcal{F}$-to-Borel measurable mapping $\gamma: \Omega \rightarrow S$ that satisfies $\gamma(\omega) \in \Phi(\omega)$ for every $\omega \in \Omega$. Since ( $S, d_{m}$ ) is locally compact, it is in particular Polish, so by the measurable selection theorem of Kuratowski and Ryll-Nardzewski [KRN65] (see also [Wag77] or [Sri98, Chapter 5.2]) it suffices to check that $\Phi(\omega)$ is nonempty and closed for every $\omega \in \Omega$, and that $\{\omega \in \Omega: E \cap \Phi(\omega)=\varnothing\} \in \mathcal{F}$ for every closed $E \subseteq S$. Since $S$ is locally compact, every closed subset of $S$ is a countable union of compact subsets, so it suffices to check the latter requirement for compact subsets of $S$, i.e., to show that $\{\omega \in \Omega: K \cap \Phi(\omega)=\varnothing\} \in \mathcal{F}$ for every compact $K \subseteq S$.

Fix $\omega \in \Omega$. If $\Gamma(\omega)=\varnothing$ then $\Phi(\omega)=S$ is closed (since $S$ is locally compact) and nonempty by assumption. If $\Gamma(\omega) \neq \varnothing$ then the continuity of the mapping $s \mapsto d_{m}(s, \Gamma(\omega))$ on $S$ implies that $\Phi(\omega)$ is closed. Moreover, in this case since $\Gamma(\omega)$ is bounded and $S$ is locally compact, the continuous mapping $s \mapsto d_{m}(s, \Gamma(\omega))$ attains its minimum on $S$, so that $\Phi(\omega) \neq \varnothing$.

It therefore remains to check that $\{\omega \in \Omega: K \cap \Phi(\omega)=\varnothing\} \in \mathcal{F}$ for every nonempty compact $K \subsetneq S$. Fixing such a $K$, since $S$ is locally compact and hence separable, there exist $\left\{\kappa_{i}\right\}_{i=1}^{\infty} \subseteq K$ and $\left\{\sigma_{j}\right\}_{j=1}^{\infty} \subseteq S$ that are dense in $K$ and $S$, respectively. Denote $\mathcal{G}_{\Gamma}=\{\omega \in \Omega: \Gamma(\omega) \neq \varnothing\}$. Then $\mathcal{G}_{\Gamma} \in \mathcal{F}$, because $\Gamma$ is strongly measurable. Observe that the following identity holds:

$$
\begin{align*}
\{\omega \in \Omega: K \cap \Phi(\omega)=\varnothing\} & =\left\{\omega \in \mathcal{G}_{\Gamma}: \forall \kappa \in K, d_{m}(\kappa, \Gamma(\omega))>d_{m}(S, \Gamma(\omega))\right\} \\
& =\bigcup_{m=1}^{\infty} \bigcap_{i=1}^{\infty} \bigcup_{j=1}^{\infty}\left\{\omega \in \mathcal{G}_{\Gamma}: d_{m}\left(\kappa_{i}, \Gamma(\omega)\right)>d_{m}\left(\sigma_{j}, \Gamma(\omega)\right)+\frac{1}{m}\right\} . \tag{232}
\end{align*}
$$

The verification of (232) proceeds as follows. Since $\Phi(\omega) \neq \varnothing$ for every $\omega \in \Omega$ and $K \neq \varnothing$, if $K \cap \Phi(\omega)=\varnothing$ then $\omega \in \mathcal{G}_{\Gamma}$ (otherwise $\Phi(\omega)=S$ ). This explains the first equality (232). For the second equality in (232), note that since $\Gamma(\omega)$ is bounded and $K$ is compact, $\inf _{\kappa \in K} d_{m}(\kappa, \Gamma(\omega))$ is attained. Therefore the second set in (232) is equal to $A=\left\{\omega \in \mathcal{G}_{\Gamma}: d_{m}(K, \Gamma(\omega))>d_{m}(S, \Gamma(\omega))\right\}$. If $\omega \in A$, then there is $m \in \mathbb{N}$ such that $d_{m}(K, \Gamma(\omega))>d_{m}(S, \Gamma(\omega))+2 / m$, implying in particular that $d_{m}\left(\kappa_{i}, \Gamma(\omega)\right)>d_{m}(S, \Gamma(\omega))+2 / m$ for every $i \in \mathbb{N}$. As $\left\{\sigma_{j}\right\}_{j=1}^{\infty}$ is dense in $S$, for every $i \in \mathbb{N}$ there is $j \in \mathbb{N}$ such that $d_{m}\left(\kappa_{i}, \Gamma(\omega)\right)>d_{m}\left(\sigma_{j}, \Gamma(\omega)\right)+1 / m$. Hence, the second set in 232] is contained in the third set in (232). For the reverse inclusion, if $\omega$ is in third set in (232) then $d m(K, \Gamma(\omega))=\inf _{i \in \mathbb{N}} d m\left(\kappa_{i}, \Gamma(\omega)\right)>\inf _{j \in \mathbb{N}} d_{m}\left(\sigma_{j}, \Gamma(\omega)\right)=d_{m}(S, \Gamma(\omega))$.

By (232), it suffices to show that $\left\{\omega \in \mathcal{G}_{\Gamma}: d_{m}(x, \Gamma(\omega))>d_{m}(y, \Gamma(\omega))+r\right\} \in \mathcal{F}$ for every fixed $x, y \in S$ and $r>0$. For this, it suffices to show that for every $z \in M$ the mapping $\omega \mapsto d m(z, \Gamma(\omega))$ is $\mathcal{F}$-to-Borel measurable on $\mathcal{G}_{\Gamma}$. Since $\mathcal{G}_{\Gamma} \in \mathcal{F}$, this is a consequence of the strong measurability of $\Gamma$, because for every $t \geqslant 0$ we have $\left\{\omega \in \mathcal{G}_{\Gamma}: d_{m}(z, \Gamma(\omega))>t\right\}=\bigcup_{k=1}^{\infty} \mathcal{G}_{\Gamma} \cap\{\omega \in \Omega: B m(z, t+1 / k) \cap \Gamma(\omega)=\varnothing\}$.

Lemma 114 is a satisfactory treatment of measurable nearest point selectors for strongly measurable set-valued mappings, though under an assumption of local compactness. We did not investigate the minimal assumptions that are required for the conclusion of Lemma 114 to hold. We will next treat the setting of standard set-valued mappings without assuming local compactness.

Let $\left(\mathscr{L}, d_{\mathscr{L}}\right)$ be a Polish metric space. Recall that a subset $A$ of $\mathscr{L}$ is said to be universally measurable if it is measurable with respect to every complete $\sigma$-finite Borel measure $\mu$ on $\mathscr{Z}$ (see e.g. [Kec95, page 155]). If $\left(m, d_{m}\right)$ is another metric space and $\Omega \subseteq \mathscr{L}$ is Borel, then a mapping $\psi: \Omega \rightarrow m$ is said to be universally measurable if $\psi^{-1}(E)$ is a universally measurable subset of $\Omega$ for every Borel subset $E$ of $m$. Finally, recall that $A \subseteq M$ is said to be analytic if it is an image under a continuous mapping of a Borel subset of a Polish metric space (see e.g. [Kec95, Chapter 14] or [Jec03, Chapter 11]). By Lusin's theorem [Luz17, Lus72] (see also e.g. [Kec95, Theorem 21.10]), analytic subsets of Polish metric spaces are universally measurable.

Lemma 115. Let $\left(m, d_{m}\right)$ and $\left(\mathscr{L}, d_{\mathcal{L}}\right)$ be Polish metric spaces and fix a Borel subset $\Omega \subseteq \mathcal{L}$. Fix also $\Delta>0$ such that $\operatorname{diam}(m) \geqslant \Delta$. Suppose that $\Gamma: \Omega \rightarrow 2^{m}$ satisfies the following two properties.
(1) For every $\omega \in \Omega$ such that $\Gamma(\omega) \neq \varnothing$ we have $\operatorname{diam} m(\Gamma(\omega))<\Delta$.
(2) For every $x \in m$ and $t \in \mathbb{R}$ the set $\left\{\omega \in \Omega: \Gamma(\omega) \neq \varnothing \wedge d_{m}(x, \Gamma(\omega))>t\right\}$ is analytic.

Then, for every closed $\varnothing \neq S \subseteq m$ there is a universally measurable mapping $\gamma: \Omega \rightarrow S$ such that

$$
\forall(\omega, x) \in \Omega \times m, \quad x \in \Gamma(\omega) \Longrightarrow d_{m}(x, \gamma(\omega)) \leqslant d_{m}(x, S)+\Delta .
$$

Proof. For every $\omega \in \Omega$, define a subset $\Psi(\omega) \subseteq S$ as follows.

$$
\Psi(\omega) \stackrel{\text { def }}{=}\left\{\begin{array}{cl}
\cap_{x \in}\left\{s \in S: d_{m}(x, s) \leqslant 2 d_{m}(x, \Gamma(\omega))+d_{m}(x, S)+\Delta\right\} & \text { if } \Gamma(\omega) \neq \varnothing  \tag{233}\\
S & \text { if } \Gamma(\omega)=\varnothing
\end{array}\right.
$$

We will show that there exists a universally measurable mapping $\gamma: \Omega \rightarrow S$ such that $\gamma(\omega) \in \Psi(\omega)$ for every $\omega \in \Omega$. Since $S$ is a closed subset of $M$, it is Polish. Hence, by the Kuratowski-Ryll-Nardzewski measurable selection theorem [KRN65], it suffices to prove that $\Psi(\omega)$ is nonempty and closed for every $\omega \in \Omega$, and that $\Psi^{-}(E)=\{\omega \in \Omega: E \cap \Psi(\omega) \neq \varnothing\}$ is universally measurable for every closed $E \subseteq S$.

By design, $\Psi(\omega)=S$ is nonempty and closed if $\Gamma(\omega)=\varnothing$. So, fix $\omega \in \Omega$ such that $\Gamma(\omega) \neq \varnothing$. Then $\Psi(\omega)$ is closed because if $\left\{s_{k}\right\}_{k=1}^{\infty} \subseteq \Psi(\omega)$ and $s \in M$ satisfy $\lim _{k \rightarrow \infty} d_{m}\left(s_{k}, s\right)=0$, then for every $k \in \mathbb{N}$ and $x \in m$, since $s_{k} \in \Psi(\omega)$ we have $d_{m}\left(s_{k}, x\right) \leqslant 2 d_{m}(x, \Gamma(\omega))+d_{m}(x, S)+\Delta$. Hence, by continuity also $d_{m}(s, x) \leqslant 2 d_{m}(x, \Gamma(\omega))+d_{m}(x, S)+\Delta$ for every $x \in m$, i.e., $s \in \Psi(\omega)$.

We will next check that $\Psi(\omega) \neq \varnothing$ for every $\omega \in \Omega$ such that $\Gamma(\omega) \neq \varnothing$. Denote $\varepsilon_{\omega}=\Delta-\operatorname{diam}_{m}(\Gamma(\omega))$. By assumption (1) of Lemma 115 we have $\varepsilon_{\omega}>0$, so we may choose $s_{\omega} \in S$ and $y_{\omega} \in \Gamma(\omega)$ that satisfy $d_{m}\left(y_{\omega}, s_{\omega}\right) \leqslant d_{m}(\Gamma(\omega), S)+\varepsilon_{\omega}$. We claim that $s_{\omega} \in \Psi(\omega)$. Indeed, for every $x \in T$ and $z \in \Gamma(\omega)$ we have

$$
\left.\begin{array}{rl}
d_{m}\left(x, s_{\omega}\right) \leqslant d_{m}(x, z)+d_{m}(z, & \left.y_{\omega}\right)+d_{m}\left(y_{\omega}, s_{\omega}\right) \tag{234}
\end{array} \leqslant d_{m}(x, z)+\operatorname{diam} m(\Gamma(\omega))+d_{m}(\Gamma(\omega), S)+\varepsilon_{\omega}\right)
$$

where in the penultimate step of (234) we used the fact that $d_{m}(\Gamma(\omega), S) \leqslant d_{m}(z, S)$, since $z \in \Gamma(\omega)$, and in the final step of (234) we used the fact that the mapping $p \mapsto d_{m}(p, S)$ is 1-Lipschitz on $m$. Since (234) holds for every $z \in \Gamma(\omega)$, it follows that $d_{m}\left(x, s_{\omega}\right) \leqslant 2 d_{m}(x, \Gamma(\omega))+d_{m}(x, S)+\Delta$. Because this holds for every $x \in M$, it follows that $s_{\omega} \in \Psi(\omega)$.

Having checked that $\Psi$ takes values in closed and nonempty subsets of $S$, it remains to show that $\Psi^{-}(E)$ is universally measurable for every closed $E \subseteq S$. To this end, since $M$ is separable, we may fix from now on a sequence $\left\{x_{j}\right\}_{j=1}^{\infty}$ that is dense in $m$. Note that by the case $t=0$ of assumption (2p) of Lemma 115 , for every $j \in \mathbb{N}$ the following set is analytic.

$$
\left\{\omega \in \Omega: \Gamma(\omega) \neq \varnothing \wedge d_{m}\left(x_{j}, \Gamma(\omega)\right)>0\right\}=\left\{\omega \in \Omega: \Gamma(\omega) \neq \varnothing \wedge x_{j} \notin \overline{\Gamma(\omega)}\right\}
$$

Countable unions and intersections of analytic sets are analytic (see e.g. [Kec95, Proposition 14.4]), so we deduce that the following set is analytic.

$$
\begin{align*}
\bigcup_{j=1}^{\infty}\left\{\omega \in \Omega: \Gamma(\omega) \neq \varnothing \wedge x_{j} \notin \overline{\Gamma(\omega)}\right\} & =\left\{\omega \in \Omega: \Gamma(\omega) \neq \varnothing \wedge\left\{x_{j}\right\}_{j=1}^{\infty} \notin \overline{\Gamma(\omega)}\right\}  \tag{235}\\
& =\{\omega \in \Omega: \Gamma(\omega) \neq \varnothing\},
\end{align*}
$$

where for the final step of (235) observe that, since $\left\{x_{j}\right\}_{j=1}^{\infty}$ is dense in $M$, if $\left\{x_{j}\right\}_{j=1}^{\infty}$ were a subset of $\overline{\Gamma(\omega)}$ then it would follow that $\Gamma(\omega)$ is dense in $M$. This would imply that $\operatorname{diam} m(\Gamma(\omega))=\operatorname{diam}(M) \geqslant \Delta$, in contradiction to assumption (1) of Lemma 115. We have thus checked that the set $\mathcal{G}_{\Gamma}=\{\omega \in \Omega: \Gamma(\omega) \neq \varnothing\}$ is analytic, and hence by Lusin's theorem Luz17, Lus72] it is universally measurable. Now,

$$
\Psi^{-}(E) \stackrel{233}{-}\left(\Omega \backslash \mathcal{G}_{\Gamma}\right) \cup\left\{\omega \in \mathcal{G}_{\Gamma}: \exists s \in E \forall x \in M, d_{m}(x, s) \leqslant 2 d_{m}(x, \Gamma(\omega))+d_{m}(x, S)+\Delta\right\}
$$

Hence, it remains to prove that the following set is universally measurable.

$$
\begin{align*}
\left\{\omega \in \mathcal{G}_{\Gamma}: \exists s \in E\right. & \left.\forall x \in m, d_{m}(x, s) \leqslant 2 d_{m}(x, \Gamma(\omega))+d_{m}(x, S)+\Delta\right\} \\
& =\left\{\omega \in \mathcal{G}_{\Gamma}: \exists s \in E \forall j \in \mathbb{N}, d_{m}\left(x_{j}, s\right) \leqslant 2 d_{m}\left(x_{j}, \Gamma(\omega)\right)+d_{m}\left(x_{j}, S\right)+\Delta\right\}, \tag{236}
\end{align*}
$$

where we used the fact that $\left\{x_{j}\right\}_{j=1}^{\infty}$ is dense in $M$.
Consider the following subset $\mathcal{C}$ of $\Omega \times E$.

$$
C \stackrel{\text { def }}{=}\left\{(\omega, s) \in \mathcal{G}_{\Gamma} \times E: \forall j \in \mathbb{N}, d_{m}\left(x_{j}, s\right) \leqslant 2 d_{m}\left(x_{j}, \Gamma(\omega)\right)+d_{m}\left(x_{j}, S\right)+\Delta\right\} .
$$

The set in 236 is $\pi_{1}(C)$, where $\pi_{1}: \Omega \times E \rightarrow \Omega$ is the projection to the first coordinate, i.e., $\pi_{1}(\omega, s)=\omega$ for every $(\omega, s) \in \Omega \times E$. Since continuous images and preimages of analytic sets are analytic (see e.g. Kec95, Proposition 14.4]), by another application of Lusin's theorem it suffices to show that $C$ is analytic. We already proved that $\mathcal{G}_{\Gamma} \subseteq \Omega$ is analytic, so there is a Borel subset $L$ of a Polish space $\mathscr{V}$ and a continuous mapping $\phi: L \rightarrow \Omega$ such that $\phi(L)=\mathcal{G}_{\Gamma}$. Denoting the identity mapping on $E$ by $\operatorname{ld}_{E}: E \rightarrow E$, since $\phi$ maps $L$ onto $\mathcal{G}_{\Gamma}$, the set $C$ is the image under the continuous mapping $\phi \times \operatorname{ld}_{E}$ of the following subset of $Y \times E$.

$$
\begin{aligned}
\left\{(y, s) \in L \times E: \forall j \in \mathbb{N}, d_{m}\left(x_{j}, s\right)\right. & \left.\leqslant 2 d_{m}\left(x_{j}, \Gamma(\phi(y))\right)+d_{m}\left(x_{j}, S\right)+\Delta\right\} \\
& =\bigcap_{j=1}^{\infty}\left\{(y, s) \in L \times E: d_{m}\left(x_{j}, s\right) \leqslant 2 d_{m}\left(x_{j}, \Gamma(\phi(y))\right)+d_{m}\left(x_{j}, S\right)+\Delta\right\}
\end{aligned}
$$

Hence, since continuous images and countable intersections of analytic sets are analytic, by yet another application of Lusin's theorem we see that it suffices to show that for every fixed $x \in M$ the following set is analytic, where for every $q \in \mathbb{Q}$ we denote $A_{q}=\left\{(y, s) \in L \times E: q<d_{m}(x, s)\right\}=L \times\{s \in E: q<d m(x, s)\}$.

$$
\begin{aligned}
& \left\{(y, s) \in L \times E: d_{m}(x, s) \leqslant 2 d_{m}(x, \Gamma(\phi(y)))+d_{m}(x, S)+\Delta\right\} \\
& =\bigcap_{q \in \mathbb{Q}}\left(\left((L \times E) \backslash A_{q}\right) \cup\left(A_{q} \cap\left\{(y, s) \in L \times E: 2 d_{m}(x, \Gamma(\phi(y)))>q-d_{m}(x, S)-\Delta\right\}\right)\right),
\end{aligned}
$$

Since $A_{q}$ is Borel for all $q \in \mathbb{Q}$, it suffices to show that the following set is analytic for every $t \in \mathbb{R}$ :

$$
\left\{(y, s) \in L \times E: d_{m}(x, \Gamma(\phi(y)))>t\right\}=\phi^{-1}\left(\left\{\omega \in \mathcal{G}_{\Gamma}: d_{m}(x, \Gamma(\omega))>t\right\}\right) \times E .
$$

Since a preimage under a continuous mapping of an analytic set is analytic, the above set is indeed analytic due to assumption (2) of Lemma 115 and the fact that $E$ is closed.

Remark 116. The proof of Lemma 115 used the assumption $\operatorname{diam}(M) \geqslant \Delta$ only to deduce that the set $\mathcal{G}_{\Gamma}=\{\omega \in \Omega: \Gamma(\omega) \neq \varnothing\}$ is analytic from (the case $t=0$ of) assumption 2) of Lemma 115. Hence, if we add the assumption that $\mathcal{G}_{\Gamma}$ is analytic to Lemma 115 , then we can drop the restriction diam $(M) \geqslant \Delta$ altogether. Alternatively, recalling equation (235) and the paragraph immediately after it, for the above proof of Lemma 115 to go through it suffices to assume that $\Gamma(\omega)$ is not dense in $M$ for any $\omega \in \Omega$.

Recalling Definition 112, Lemma 115 and Remark 116 imply the following corollary. Indeed, by Remark 116 we know that we can drop the assumption $\operatorname{diam}(M) \geqslant \Delta$ of Lemma 115 , and when $\Gamma$ is a standard set-valued mapping the sets that appears in assumption 20) of Lemma 115 are Borel.

Corollary 117. Fix $\Delta>0$. Let $\left(M, d_{m}\right)$ and $\left(\mathscr{L}, d_{\mathcal{L}}\right)$ be Polish metric spaces and fix a Borel subset $\Omega \subseteq \mathscr{L}$. Suppose that $\Gamma: \Omega \rightarrow 2^{m}$ is a standard set-valued mapping such that $\operatorname{diam} m(\Gamma(\omega))<\Delta$ for every $\omega \in \mathcal{G} \Gamma$. Then for every closed $\varnothing \neq S \subseteq m$ there exists a universally measurable mapping $\gamma: \Omega \rightarrow S$ that satisfies

$$
\forall(\omega, x) \in \Omega \times m, \quad x \in \Gamma(\omega) \Longrightarrow d m(x, \gamma(\omega)) \leqslant d m(x, S)+\Delta
$$

3.3. Measurability of iterative ball partitioning. The following set-valued mapping is a building block of much of the literature on random partitions, including the present investigation. Fix a metric space $(m, d m)$ and $k \in \mathbb{N}$. Define a set-valued mapping $\Gamma: m^{k} \times[0, \infty)^{k} \rightarrow 2^{m}$ by

$$
\begin{equation*}
\forall(\vec{x}, \vec{r})=\left(x_{1}, \ldots, x_{k}, r_{1}, \ldots, r_{k}\right) \in m^{k} \times[0, \infty)^{k}, \quad \Gamma(\vec{x}, \vec{r}) \stackrel{\text { def }}{=} B m\left(x_{k}, r_{k}\right) \backslash \bigcup_{j=1}^{k-1} B m\left(x_{j}, r_{j}\right) \tag{237}
\end{equation*}
$$

We can think of $\Gamma$ as a random subset of $M$ if we are given a probability measure Prob on $m^{k} \times[0, \infty)^{k}$. The measure Prob can encode the geometry of $(m, d m)$; for example, if $(m, d m)$ is a complete doubling metric space, then in [LN05] this measure arises from a doubling measure on $M$ (see [VK87, LS98]). The measure Prob can also have a "smoothing effect" through the randomness of the radii (see e.g. Bar99, CKR05, FRT04, LN05, MN07, NT10, ABN11, NT12]; choosing a suitable distribution over the random radii is sometimes an important and quite delicate matter, but this intricacy will not arise in the present work. For finite dimensional normed spaces, a random subset as in (237) was used in [CG ${ }^{+} 98$ KMS98]. Note that given $\Delta>0$, if the measure Prob is supported on the set of those $(\vec{x}, \vec{r}) \in m^{k} \times[0, \infty)^{k}$ for which $r_{k} \leqslant \Delta / 2$, then the mapping $\Gamma$ takes values in subsets of $M$ of diameter at most $\Delta$.

While the definition 237) is very simple and natural, in order to use it in the ensuing reasoning we need to know that it satisfies certain measurability requirements. Note first that the set-valued mapping $\Gamma$ in 237) automatically has the following basic measurability property: For every fixed $y \in M$ the set $\left\{(\vec{x}, \vec{r}) \in M^{k} \times[0, \infty)^{k}: y \in \Gamma(\vec{x}, \vec{r})\right\}$ is Borel. Indeed, by definition we have

$$
\begin{aligned}
&\left\{(\vec{x}, \vec{r}) \in m^{k} \times[0, \infty)^{k}: y \in \Gamma(\vec{x}, \vec{r})\right\} \\
&=\bigcap_{j=1}^{k-1}\left\{(\vec{x}, \vec{r}) \in m^{k} \times[0, \infty)^{k}: d_{m}\left(y, x_{j}\right)>r_{j}\right\} \cap\left\{(\vec{x}, \vec{r}) \in m^{k} \times[0, \infty)^{k}: d_{m}\left(y, x_{k}\right) \leqslant r_{k}\right\} .
\end{aligned}
$$

In other words, the indicator mapping $(\vec{x}, \vec{r}) \mapsto \mathbf{1}_{\Gamma(\vec{x}, \vec{r})}(y)$ is Borel measurable for every fixed $y \in M$.
Lemma 118. Fix $k \in \mathbb{N}$. Suppose that $\left(M, d_{m}\right)$ is a Polish metric space. Let $\Gamma: m^{k} \times[0, \infty)^{k} \rightarrow 2^{m}$ be given in 237). Then $\Gamma^{-}(S)=\left\{(\vec{x}, \vec{r}) \in m^{k} \times[0, \infty)^{k}: S \cap \Gamma(\vec{x}, \vec{r}) \neq \varnothing\right\}$ is analytic for every analytic subset $S \subseteq M$. Consequently, for every complete $\sigma$-finite Borel measure $\mu$ on $m^{k} \times[0, \infty)^{k}$, if $\mathcal{F}_{\mu}$ denotes the $\sigma$-algebra of $\mu$-measurable subsets of $m^{k} \times[0, \infty)^{k}$, then $\Gamma$ is a strongly measurable set-valued mapping from the measurable space $\left(m^{k} \times[0, \infty)^{k}, \mathcal{F}_{\mu}\right)$ to $2^{m}$.

Proof. Since $S$ is analytic, there exists a Borel subset $T$ of a Polish metric space $\mathscr{L}$ and a continuous mapping $\psi: T \rightarrow M$ such that $\psi(T)=S$. Consider the following Borel subset $\mathscr{B}$ of the Polish space $m^{k} \times[0, \infty)^{k} \times \mathscr{L}(\mathscr{B}$ is Borel because it is defined using finitely many continuous inequalities).

$$
\mathcal{B} \stackrel{\operatorname{def}}{=}\left\{(\vec{x}, \vec{r}, t) \in m^{k} \times[0, \infty)^{k} \times T: d_{m}\left(\psi(t), x_{k}\right) \leqslant r_{k} \wedge \forall j \in\{1, \ldots, k-1\}, d_{m}\left(\psi(t), x_{j}\right)>r_{j}\right\} .
$$

Then $\Gamma^{-}(S)=\pi(\mathscr{B})$, where $\pi: m^{k} \times[0, \infty)^{k} \times \mathscr{L} \rightarrow m^{k} \times[0, \infty)^{k}$ is the projection onto the first two coordinates, i.e., $\pi(\vec{x}, \vec{r}, z)=(\vec{x}, \vec{r})$ for every $(\vec{x}, \vec{r}, z) \in M^{k} \times[0, \infty)^{k} \times \mathcal{L}$. Since $\pi$ is continuous, it follows that $\Gamma^{-}(S)$ is analytic. By Lusin's theorem [Luz17, Lus72], it follows that $\Gamma^{-}(S)$ is universally measurable. In particular, if $\mu$ is a complete $\sigma$-finite Borel measure on $\Pi^{k} \times[0, \infty)^{k}$ and $\mathcal{F}_{\mu}$ is the $\sigma$-algebra of $\mu$-measurable subsets of $M^{k} \times[0, \infty)^{k}$, then $\Gamma^{-}(E) \in \mathcal{F}_{\mu}$ for every closed subset $E \subseteq M$. Recalling (92), this means that $\Gamma$ is a strongly measurable set-valued mapping from the measurable space $\left(m^{k} \times[0, \infty)^{k}, \mathcal{F}_{\mu}\right)$ to $2^{m}$.

Lemma 119 below contains additional Borel measurability assertions that will be used later. Its assumptions are satisfied, for example, when $m$ is a separable normed space, which is the case of interest here. We did not investigate the maximal generality under which the conclusion of Lemma 119 holds.

In what follows, given a metric space $(m, d m)$, for every $x \in m$ and $r>0$ the open ball of radius $r$ centered at $x$ is denoted $B_{m}^{\circ}(x, r)=\left\{y \in m: d_{m}(x, y)<r\right\}$.
Lemma 119. Suppose that $\left(m, d_{m}\right)$ is a separable metric space such that

$$
\begin{equation*}
\forall(x, r) \in M \times(0, \infty), \quad B m(x, r)=\overline{B_{m}^{\circ}(x, r)} . \tag{238}
\end{equation*}
$$

Fix $k \in \mathbb{N}$ and let $\Gamma: m^{k} \times(0, \infty)^{k} \rightarrow 2^{m}$ be given in (237). Then the following set is Borel measurable.

$$
\mathcal{G}_{\Gamma}=\left\{(\vec{x}, \vec{r}) \in M^{k} \times(0, \infty)^{k}: \Gamma(\vec{x}, \vec{r}) \neq \varnothing\right\} .
$$

Also, for each $y \in T$ the mapping from $\mathcal{G}_{\Gamma}$ to $\mathbb{R}$ that is given by $(\vec{x}, \vec{r}) \mapsto d_{m}(y, \Gamma(x, r))$ is Borel measurable.
Proof. Let $\mathfrak{D} \subseteq m$ be a countable dense subset of $m$. The assumption (238) implies that $\mathfrak{D} \cap \Gamma(\vec{x}, \vec{r})$ is dense in $\Gamma(\vec{x}, \vec{r})$ for every $(\vec{x}, \vec{r}) \in m^{k} \times(0, \infty)^{k}$. This is straightforward to check as follows. Fix $y \in \Gamma(\vec{x}, \vec{r})$ and $\delta>0$. We need to find $q \in \mathfrak{D} \cap \Gamma(\vec{x}, \vec{r})$ with $d_{m}(q, y)<\delta$. Recalling (237), since $y \in \Gamma(\vec{x}, \vec{r})$ we know that $d_{m}\left(y, x_{k}\right) \leqslant r_{k}$, and also $d_{m}\left(y, x_{j}\right)>r_{j}$ for every $j \in\{1, \ldots, k-1\}$, i.e., $\eta>0$ where

$$
\eta \stackrel{\text { def }}{=} \min \left\{\delta, d_{m}\left(y, x_{1}\right)-r_{1}, \ldots, d_{m}\left(y, x_{k-1}\right)-r_{k-1}\right\} .
$$

By (238) there is $z \in B_{m}^{\circ}\left(x_{k}, r_{k}\right)$ with $d_{m}(z, y)<\eta / 2$. Denote

$$
\rho \stackrel{\text { def }}{=} \min \left\{r_{k}-d_{m}\left(z, x_{k}\right), \frac{1}{2} \eta\right\} .
$$

Then $\rho>0$, so the density of $\mathfrak{D}$ in $m$ implies that there is $q \in \mathfrak{D}$ with $d_{m}(q, z)<\rho$. Consequently,

$$
d_{m}(q, y) \leqslant d_{m}(q, z)+d_{m}(z, y)<\rho+\frac{\eta}{2} \leqslant \delta
$$

It remains to observe that $q \in \Gamma(\vec{x}, \vec{r})$, because $d_{m}\left(q, x_{k}\right) \leqslant d_{m}(q, z)+d_{m}\left(z, x_{k}\right)<\rho+d_{m}\left(z, x_{k}\right) \leqslant r_{k}$ and also for every $j \in\{1, \ldots, k-1\}$ we have

$$
d_{m}\left(q, x_{j}\right) \geqslant d_{m}\left(y, x_{j}\right)-d_{m}(y, z)-d_{m}(z, q)>d_{m}\left(y, x_{j}\right)-\frac{\eta}{2}-\rho \geqslant d_{m}\left(y, x_{j}\right)-\eta \geqslant r_{j}
$$

For every $(\vec{x}, \vec{r}) \in T^{k} \times(0, \infty)^{k}$, we have $\Gamma(\vec{x}, \vec{r}) \neq \varnothing$ if and only if $\mathfrak{D} \cap \Gamma(\vec{x}, \vec{r}) \neq \varnothing$. Consequently,

$$
\mathcal{G}_{\Gamma}=\left\{(\vec{x}, \vec{r}) \in m^{k} \times(0, \infty)^{k}: \Gamma(\vec{x}, \vec{r}) \neq \varnothing\right\}=\bigcup_{q \in \mathfrak{D}}\left\{(\vec{x}, \vec{r}) \in m^{k} \times(0, \infty)^{k}: q \in \Gamma(\vec{x}, \vec{r})\right\} .
$$

Since $\mathfrak{D}$ is countable and we already checked in the paragraph immediately preceding Lemma 118 that $\left\{(\vec{x}, \vec{r}) \in M^{k} \times(0, \infty)^{k}: y \in \Gamma(\vec{x}, \vec{r})\right\}$ is Borel measurable for every $y \in M$, we get that $\mathcal{G}_{\Gamma}$ is Borel measurable.

Next, $d_{m}(y, \Gamma(\vec{x}, \vec{r}))=d_{m}(y, \mathfrak{D} \cap \Gamma(\vec{x}, \vec{r}))$ for every $(\vec{x}, \vec{r}) \in \mathcal{G}_{\Gamma}$ and $y \in m$. So, for every $t>0$ we have

$$
\left\{(\vec{x}, \vec{r}) \in \mathcal{G}_{\Gamma}: d m(y, \Gamma(\vec{x}, \vec{r}))<t\right\}=\bigcup_{q \in \mathfrak{D} \cap B_{m}^{o}(y, t)}\left\{(\vec{x}, \vec{r}) \in m^{k} \times(0, \infty)^{k}: q \in \Gamma(\vec{x}, \vec{r})\right\} .
$$

It follows that $\left\{(\vec{x}, \vec{r}) \in \mathcal{G}_{\Gamma}: d_{m}(y, \Gamma(\vec{x}, \vec{r}))<t\right\}$ is Borel measurable for every $t \in \mathbb{R}$.
Corollary 120 below follows directly from the definition of a standard set-valued mapping due to Lemma 119 and the discussion in the paragraph immediately preceding Lemma 118 .
Corollary 120. Let ( $M, d_{m}$ ) be a Polish metric space satisfying (238). Then, for every $k \in \mathbb{N}$ the set-valued mapping $\Gamma: m^{k} \times(0, \infty)^{k} \rightarrow 2^{m}$ in (237) is standard.

## 4. UPPER BOUNDS ON RANDOM PARTITIONS

In this section, we will prove the existence of random partitions with the separation and padding properties that were stated in the Introduction.
4.1. Proof of Theorem 74 and the upper bound on $\operatorname{PAD}_{\delta}(\mathbf{X})$ in Theorem 68 . Theorem 121 below asserts that every normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ admits a random partition that simultaneously has desirable padding and separation properties. In the literature, such properties are obtained for different random partitions: Separating partitions of normed spaces use iterative ball partitioning with deterministic radii, while padded partitions also rely on randomizing the radii. At present, we do not have in mind an application in which good padding and separation properties are needed simultaneously for the same random partition, so it is worthwhile to note this feature for potential future use but in what follows we will use Theorem 121 to obtain two standalone conclusions that yield upper bounds on the moduli of padded and separated decomposability (in fact, the separation profile of Theorem 74 ).

Theorem 121. Fix $n \in \mathbb{N}$ and a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$. For every $\Delta \in(0, \infty)$ there exists $a \Delta$-bounded random partition $\mathcal{P}_{\Delta}$ of $\mathbf{X}$ such that for every $x, y \in \mathbb{R}^{n}$ and every $\delta \in(0,1)$ we have

$$
\begin{equation*}
\operatorname{Prob}\left[\mathcal{P}_{\Delta}(x) \neq \mathcal{P}_{\Delta}(y)\right]=\min \left\{1, \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{(x-y)^{\perp}}(B \mathbf{X})\right)}{\Delta \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}\|x-y\|_{\ell_{2}^{n}}\right\} \tag{239}
\end{equation*}
$$

and,

$$
\operatorname{Prob}\left[\mathcal{P}_{\Delta}(x) \supseteq \frac{1-\sqrt[n]{\delta}}{1+\sqrt[n]{\delta}} \cdot \frac{\Delta}{2} B_{\mathbf{X}}\right]=\delta
$$

By the conventions of Remark 61 , the $\Delta$-boundedness of Theorem 121 is with respect to the norm $\|\cdot\|_{\mathbf{x}}$, i.e., the clusters of the random partition $\mathcal{P}_{\Delta}$ have $\mathbf{X}$-diameter at most $\Delta$. By the definitions in Section 1.7.1, the notion of random partition implies that each of the clusters of $\mathcal{P}_{\Delta}$ is strongly measurable, but we will see that they are also standard (recall Definition 112 .
Remark 122. For every $M>0$, consider the metric space $L_{1}^{\leqslant M}=\left(L_{1}, d_{M}\right)$ that is given by

$$
\forall f, \in L_{1}, \quad d_{M}(f, g) \stackrel{\text { def }}{=} \min \left\{M,\|f-g\|_{L_{1}}\right\}
$$

A useful property MN15, Lemma 5.4] of this truncated $L_{1}$ metric is $c_{L_{1}}\left(L_{1}^{\leqslant M}\right) \lesssim 1$, i.e., $L_{1}^{\leqslant M}$ embeds back into $L_{1}$ with bi-Lipschitz distortion $O(1)$. Theorem 121 gives a different proof of this since if $\mathbf{X}=\ell_{\infty}^{n}$, then by (38) the right hand side of (239) is equal to $\min \left\{2 \Delta,\|x-y\|_{1}\right\} /(2 \Delta)$. At the same time, if $\mathcal{P}_{\Delta}^{\omega}=\left\{\Gamma_{\Delta}^{k}(\omega)\right\}_{k=1}^{\infty}$, then the left hand side of 239 ) embeds isometrically into an $L_{1}(\mu)$ space via the embedding

$$
\left(f \in L_{1}\right) \mapsto\left(\omega \mapsto\left(\mathbf{1}_{\Gamma^{k}(\omega)}(f)\right)_{k=1}^{\infty}\right) \in L_{1}\left(\text { Prob; } \ell_{1}\right)
$$

By (30), the right hand side of (239) equals $\min \left\{\Delta,\|x-y\|_{\Pi^{*} x}\right\} / \Delta$. However, the class of finite dimensional normed spaces whose unit ball is a polar projection body coincides with Bol69 those finite dimensional normed spaces that embed isometrically into $L_{1}$, so this does not give a new embedding result.

We will first describe the construction that leads to the random partition whose existence is asserted in Theorem 121. This construction is a generalization of the construction that appears in the proof Lemma 3.16 of [LN05], which itself combines a coloring argument with a generalization of the iterated ball partitioning technique that was used in the Euclidean setting in [CCG ${ }^{+} 98$, KMS98].

In the rest of this section we will work under the assumptions and notation of Theorem 121 , Let $\Lambda \subseteq \mathbb{R}^{n}$ be a lattice such that $\left\{z+B_{\mathbf{X}}\right\}_{z \in \Lambda}$ have pairwise disjoint interiors (equivalently, $\left\|z-z^{\prime}\right\|_{\mathbf{X}} \geqslant 2$ for every distinct $z, z^{\prime} \in \Lambda$ ) and $\bigcup_{z \in \Lambda}\left(z+3 B_{\mathbf{X}}\right)=\mathbb{R}^{n}$ (i.e., for every $x \in \mathbb{R}^{n}$ there is $z \in \Lambda$ such that $\|x-z\|_{\mathbf{X}} \leqslant 3$ ). The existence of such a lattice follows from the work of Rogers [Rog50] (see [Zon02, Remark 6]). The constant 3 here is not the best-known (see But72, Zon02); we prefer to work with an explicit constant only for notational convenience despite the fact that its value is not important in the present context.

Denote the $\mathbf{X}$-Voronoi cell of $\Lambda$, i.e., the set of points in $\mathbb{R}^{n}$ whose closest lattice point is the origin, by

$$
\mathcal{V} \stackrel{\text { def }}{=}\left\{x \in \mathbb{R}^{n}:\|x\|_{\mathbf{X}}=\min _{z \in \Lambda}\|x-z\|_{\mathbf{X}}\right\} .
$$

Then $\mathcal{V} \subseteq 3 B_{\mathbf{X}}$ and the translates $\{z+\mathcal{V}\}_{z \in \Lambda}$ cover $\mathbb{R}^{n}$ and have pairwise disjoint interiors.

Remark 123. Our choice of the above lattice is natural since it is adapted to the intrinsic geometry of $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ and it leads to a simpler probability space in the construction below. Nevertheless, for the present purposes this choice is not crucial, and one could also work with any other lattice, including $\mathbb{Z}^{n}$. In that case, one could carry out the ensuing reasoning while adapting it to geometric characteristics of the lattice in question (its packing radius, covering radius and the diameter of its Voronoi cell, all of which are measured with respect to the metric induced by $\|\cdot\|_{\mathbf{X}}$ ). This requires several changes in the ensuing discussion, resulting in slightly more cumbersome computations that incorporate these geometric characteristics of the lattice. All of these quantities are universal constants for our choice of $\Lambda$.

Define graph $\mathrm{G}=\left(\Lambda, \mathrm{E}_{\mathrm{G}}\right)$ whose vertex set is the lattice $\Lambda$ and whose edge set $\mathrm{E}_{\mathrm{G}}$ is given by

$$
\forall w, z \in \Lambda, \quad\{w, z\} \in \mathrm{E}_{\mathrm{G}} \Longleftrightarrow w \neq z \wedge \inf _{\substack{a \in w+\mathcal{V} \\ b \in z+\mathcal{V}}}\|a-b\|_{\mathbf{X}} \leqslant 10
$$

So, if $\{w, z\} \in \mathrm{E}_{\mathrm{G}}$ and $x \in B_{\mathbf{X}}$ then there exist $u, v \in \mathcal{V}$ such that $\|(w+u)-(z+v)\|_{\mathbf{X}} \leqslant 10$ and therefore, since $\mathcal{V} \subseteq 3 B_{\mathbf{X}}$, we have $\|w-(z+x)\|_{\mathbf{X}} \leqslant\|(w+u)-(z+v)\|_{\mathbf{X}}+\|u\|_{\mathbf{X}}+\|v\|_{\mathbf{X}}+\|x\|_{\mathbf{X}} \leqslant 17$. Hence $z+B_{\mathbf{X}} \subseteq w+17 B_{\mathbf{X}}$. It follows that if $w \in \Lambda$ and $z_{1}, \ldots, z_{m} \in \Lambda$ are the distinct neighbors of $w$ in the graph $G$ then the balls $\left\{z_{i}+B_{\mathbf{X}}\right\}_{i=1}^{m}$ have disjoint interiors (since distinct elements of the lattice $\Lambda$ are at $\mathbf{X}$-distance at least 2), yet they are all contained in the ball $w+17 B_{\mathbf{X}}$. By comparing volumes, this implies that $m \leqslant 17^{n}$. In other words, the degree of the graph G is at most $17^{n}$, and therefore (by applying the greedy algorithm, see e.g. Bro41]) its chromatic number is at most $17^{n}+1 \leqslant 5^{2 n}$, i.e., there is $\chi: \Lambda \rightarrow\left\{1, \ldots, 5^{2 n}\right\}$ such that

$$
\begin{equation*}
\forall w, z \in \Lambda, \quad w \neq z \wedge \inf _{\substack{a \in w+\mathcal{V} \\ b \in z+\mathcal{V}}}\|a-b\|_{\mathbf{x}} \leqslant 10 \Longrightarrow \chi(w) \neq \chi(z) \tag{240}
\end{equation*}
$$

Consider the Polish space $\mathscr{L} \stackrel{\text { def }}{=} \mathcal{V}^{\mathbb{N}} \times\left\{1, \ldots, 5^{2 n}\right\}^{\mathbb{N}}$. In what follows, every $\omega \in \mathscr{L}$ will be written as $\omega=(\vec{x}, \vec{\gamma})$, where $\vec{x}=\left(x_{1}, x_{2}, \ldots\right) \in \mathcal{V}^{\mathbb{N}}$ and $\vec{\gamma}=\left(\gamma_{1}, \gamma_{2}, \ldots\right) \in\left\{1, \ldots, 5^{2 n}\right\}^{\mathbb{N}}$. Denote by $\mu$ the normalized Lebesgue measure on $\mathcal{V}$ and by $v$ the normalized counting measure on $\left\{1, \ldots, 5^{2 n}\right\}$, i.e., for every Lebesgue measurable $A \subseteq \mathbb{R}^{n}$ and every $F \subseteq\left\{1, \ldots, 5^{2 n}\right\}^{\mathbb{N}}$ we have

$$
\mu(A) \stackrel{\text { def }}{=} \frac{\operatorname{vol}_{n}(A \cap \mathcal{V})}{\operatorname{vol}_{n}(\mathcal{V})} \quad \text { and } \quad v(F) \stackrel{\text { def }}{=} \frac{|F|}{5^{2 n}}
$$

Henceforth, the product probability measure $\mu^{\mathbb{N}} \times v^{\mathbb{N}}$ on $\mathscr{L}$ will be denoted by Prob.
For every $k \in \mathbb{N}, z \in \Lambda$ and $(\vec{x}, \vec{\gamma}) \in \mathscr{L}$ define a subset $\Gamma^{k, z}(\vec{x}, \vec{\gamma}) \subseteq \mathbb{R}^{n}$ by

$$
\begin{align*}
& \chi(z)=\gamma_{k} \Longrightarrow \Gamma^{k, z}(\vec{x}, \vec{\gamma}) \stackrel{\text { def }}{=}\left(z+x_{k}+B_{\mathbf{X}}\right) \backslash \bigcup_{j=1}^{k-1} \bigcup_{\substack{w \in \Lambda \\
\chi(w)=\gamma_{j}}}\left(w+x_{j}+B_{\mathbf{X}}\right)  \tag{241}\\
& \chi(z) \neq \gamma_{k} \Longrightarrow \Gamma^{k, z}(\vec{x}, \vec{\gamma}) \stackrel{\text { def }}{=} \varnothing
\end{align*}
$$

Lemma 124. For every $k \in \mathbb{N}$ and $z \in \Lambda$ the set-valued mapping $\Gamma^{k, z}: \mathscr{L} \rightarrow 2^{\mathbb{R}^{n}}$ is both strongly measurable and standard (where the underlying $\sigma$-algebra on $\mathscr{L}$ is the Prob-measurable sets).
Proof. For every $\chi_{1}, \ldots, \chi_{k} \in\left\{1, \ldots, 5^{2 n}\right\}$ consider the cylinder set

$$
C\left(\chi_{1}, \ldots, \chi_{k}\right) \stackrel{\text { def }}{=}\left\{(\vec{x}, \vec{\gamma}) \in \mathscr{L}:\left(\gamma_{1}, \ldots, \gamma_{k}\right)=\left(\chi_{1}, \ldots, \chi_{k}\right)\right\} .
$$

Since $\left\{C\left(\chi_{1}, \ldots, \chi_{k}\right):\left(\chi_{1}, \ldots, \chi_{k}\right) \in\left\{1, \ldots, 5^{2 n}\right\}^{k}\right\}$ is a partition of $\mathscr{L}$ into finitely many measurable sets, it suffices to fix from now on a $k$-tuple of colors $\vec{\chi}=\left(\chi_{1}, \ldots, \chi_{k}\right) \in\left\{1, \ldots, 5^{2 n}\right\}^{k}$ and to show that the restriction of $\Gamma^{k, z}$ to $C\left(\chi_{1}, \ldots, \chi_{k}\right)$ is both strongly measurable and standard.

Observe that for each fixed $z \in \Lambda$ and $\gamma \in\left\{1, \ldots, 5^{2 n}\right\}$ there is at most one $w \in \Lambda$ that satisfies $\chi(w)=\gamma$ and $\left(z+\mathcal{V}+B_{\mathbf{X}}\right) \cap\left(w+\mathcal{V}+B_{\mathbf{X}}\right) \neq \varnothing$. Indeed, if both $w \in \Lambda$ and $w^{\prime} \in \Lambda$ satisfied these two requirements
then we would have $\chi(w)=\gamma=\chi\left(w^{\prime}\right)$ and there would exist $a, a^{\prime}, b, b^{\prime} \in \mathcal{V}$ and $u, u^{\prime}, v, v^{\prime} \in B_{\mathbf{X}}$ such that $w+a+u=z+b+v$ and $w^{\prime}+a^{\prime}+u^{\prime}=z+b^{\prime}+v^{\prime}$. Hence,

$$
\begin{aligned}
& \inf _{\substack{\alpha \in w+v \\
\beta \in w^{\prime}+\mathcal{V}}}\|\alpha-\beta\|_{\mathbf{X}} \leqslant\left\|(w+a)-\left(w^{\prime}+a^{\prime}\right)\right\|_{\mathbf{X}}=\left\|(z+b+v-u)-\left(z+b^{\prime}+v^{\prime}-u^{\prime}\right)\right\|_{\mathbf{X}} \\
& \leqslant\|b\|_{\mathbf{X}}+\left\|b^{\prime}\right\|_{\mathbf{X}}+\|v\|_{\mathbf{X}}+\left\|v^{\prime}\right\|_{\mathbf{X}}+\|u\|_{\mathbf{X}}+\left\|u^{\prime}\right\|_{\mathbf{X}} \leqslant 3+3+1+1+1+1=10,
\end{aligned}
$$

where we used the fact that $b, b^{\prime} \in \mathcal{V} \subseteq 3 B_{\mathbf{X}}$. By (240) this contradicts the fact that $\chi(w)=\chi\left(w^{\prime}\right)$.
Having checked that the above $w$ is unique, denote it by $w(\gamma, z) \in \Lambda$. If there is no $w \in \Lambda$ that satisfies $\chi(w)=\gamma$ and $\left(z+\mathcal{V}+B_{\mathbf{X}}\right) \cap\left(w+\mathcal{V}+B_{\mathbf{X}}\right) \neq \varnothing$ then let $w(\gamma, z) \in \Lambda$ be an arbitrary (but fixed) lattice point such that $\left(z+\mathcal{V}+B_{\mathbf{x}}\right) \cap\left(w(\gamma, z)+\mathcal{V}+B_{\mathbf{X}}\right)=\varnothing$. Observe that $w(\chi(z), z)=z$. Under this notation, for every $x_{1}, \ldots, x_{k} \in \mathcal{V}$ and $\gamma_{1}, \ldots, \gamma_{k-1} \in\left\{1, \ldots, 5^{2 n}\right\}$ we have

$$
\left(z+x_{k}+B_{\mathbf{X}}\right) \backslash \bigcup_{j=1}^{k-1} \bigcup_{\substack{w \in \Lambda \\ \chi(w)=\gamma_{j}}}\left(w+x_{j}+B_{\mathbf{X}}\right)=\left(w(\chi(z), z)+x_{k}+B_{\mathbf{X}}\right) \backslash \bigcup_{j=1}^{k-1}\left(w\left(\gamma_{j}, z\right)+x_{j}+B_{\mathbf{X}}\right) .
$$

Equivalently, if we denote for every $\overrightarrow{\mathrm{y}}=\left(\mathrm{y}_{1}, \ldots, \mathrm{y}_{k}\right) \in\left(\mathbb{R}^{n}\right)^{k}$,

$$
\Theta^{k}(\vec{y}) \stackrel{\operatorname{def}}{=}\left(y_{k}+B_{\mathbf{x}}\right) \backslash \bigcup_{j=1}^{k-1}\left(y_{j}+B_{\mathbf{X}}\right),
$$

then the definition (241) can be rewritten as the assertion that the restriction of $\Gamma^{k, z}$ to $C(\vec{\chi})$ is the constant function $\varnothing$ if $\chi(z) \neq \chi_{k}$, while if $\chi(z)=\chi_{k}$ then $\Gamma^{k, z}(\vec{x}, \vec{\gamma})=\Theta^{k}(w(\vec{\chi}, z)+\vec{x})$ for every $(\vec{x}, \vec{\gamma}) \in C(\vec{\chi})$, where we use the notation $w(\vec{\chi}, z)=\left(w\left(\chi_{1}, z\right), \ldots, w\left(\chi_{k}, z\right)\right) \in\left(\mathbb{R}^{n}\right)^{k}$. The desired measurability of the restriction of $\Gamma^{k, z}$ to $C(\vec{\chi})$ now follows from Lemma 118 and Corollary 120 .

Since the sets $\{z+\mathcal{V}\}_{z \in \Lambda}$ cover $\mathbb{R}^{n}$, for every rational point $q \in \mathbb{Q}^{n}$ we can fix from now on a lattice point $z_{q} \in \Lambda$ such that $q \in z_{q}+\mathcal{V}$. Define a subset $\Omega \subseteq \mathscr{Z}=\mathcal{V}^{\mathbb{N}} \times\left\{1, \ldots, 5^{2 n}\right\}^{\mathbb{N}}$ by

$$
\begin{equation*}
\Omega \stackrel{\operatorname{def}}{=} \bigcap_{m=1}^{\infty} \bigcap_{q \in \mathbb{Q}^{n}} \bigcup_{k=1}^{\infty}\left\{(\vec{x}, \vec{\gamma}) \in \mathscr{Z}: \chi\left(z_{q}\right)=\gamma_{k} \wedge\left\|\left(z_{q}+x_{k}\right)-q\right\| \mathbf{X} \leqslant \frac{1}{m}\right\} . \tag{242}
\end{equation*}
$$

We record for ease of later use the following simple properties of $\Omega$.
Lemma 125. $\Omega$ is a Borel subset of $\mathscr{Z}$ that satisfies $\operatorname{Prob}[\Omega]=1$. Furthermore, for every $(\vec{x}, \vec{\gamma}) \in \Omega$ the set $\left\{z+x_{k}:(k, z) \in \mathbb{N} \times \Lambda \wedge \chi(z)=\gamma_{k}\right\}$ is dense in $\mathbb{R}^{n}$.

Proof. The fact that $\Omega$ is Borel is evident from its definition (242). Also, if $(\vec{x}, \vec{\gamma}) \in \Omega, u \in \mathbb{R}^{n}$ and $\varepsilon \in(0,1)$, then choose $q \in \mathbb{Q}^{n}$ such that $\|u-q\|_{\mathbf{X}}<\varepsilon / 2$. Setting $m=\lceil 2 / \varepsilon\rceil \in \mathbb{N}$, it follows from (242) that there exists $k \in \mathbb{N}$ satisfying $\chi\left(z_{q}\right)=\gamma_{k}$ and $\left\|\left(z_{q}+x_{k}\right)-q\right\|_{\mathbf{X}} \leqslant 1 / m \leqslant \varepsilon / 2$. By our choice of $q$, it follows that $\left\|\left(z_{q}+x_{k}\right)-u\right\|_{\mathbf{X}}<\varepsilon$. Since this holds for every $\varepsilon \in(0,1)$, the set $\left\{z+x_{k}:(k, z) \in \mathbb{N} \times \Lambda \wedge \chi(z)=\gamma_{k}\right\}$ is dense in $\mathbb{R}^{n}$. It remains to show that $\operatorname{Prob}[\Omega]=1$. Indeed,

$$
\begin{align*}
& \operatorname{Prob}[\mathscr{L} \backslash \Omega] \stackrel{[242}{\leqslant} \\
& \sum_{m=1}^{\infty} \sum_{q \in \mathbb{Q}^{n}} \operatorname{Prob}\left[\bigcap_{k=1}^{\infty} \mathscr{Z} \backslash\left\{(\vec{x}, \vec{\gamma}) \in \mathscr{Z}: \chi\left(z_{q}\right)=\gamma_{k} \wedge\left\|\left(z_{q}+x_{k}\right)-q\right\|_{\mathbf{x}} \leqslant \frac{1}{m}\right\}\right]  \tag{243}\\
&=\sum_{m=1}^{\infty} \sum_{q \in \mathbb{Q}^{n}} \lim _{\ell \rightarrow \infty}\left(1-\frac{\operatorname{vol}_{n}\left(\left(q-z_{q}+\frac{1}{m} B_{\mathbf{X}}\right) \cap \mathcal{V}\right)}{5^{2 n} \operatorname{vol}_{n}(\mathcal{V})}\right)^{\ell}=0,
\end{align*}
$$

where for the penultimate step of (243) recall that $\operatorname{Prob}=\mu^{\mathbb{N}} \times v^{\mathbb{N}}$. For the final step of (243) one needs to check that $\operatorname{vol}_{n}\left(\left(q-z_{q}+r B_{\mathbf{X}}\right) \cap \mathcal{V}\right)=\operatorname{vol}_{n}\left(\left(q+r B_{\mathbf{X}}\right) \cap\left(z_{q}+\mathcal{V}\right)\right)>0$ for every fixed $q \in \mathbb{Q}^{n}$ and $r \in(0, \infty)$. This is so because $z_{q} \in \Lambda$ was chosen so that $q \in z_{q}+\mathcal{V}$ (and $\mathcal{V}$ is a convex body).

The following lemma introduces the random partition that will be used to prove Theorem 121 .

Lemma 126. $\mathcal{P} \stackrel{\text { def }}{=}\left\{\left.\Gamma^{k, z}\right|_{\Omega}: \Omega \rightarrow 2^{\mathbb{R}^{n}}\right\}_{(k, z) \in \mathbb{N} \times \Lambda}$ is a 2 -bounded random partition of $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$, each of whose clusters are both strongly measurable and standard set-valued mappings.

Proof. Since $\Omega$ is a Borel subset of $\mathscr{L}$, for each $(k, z) \in \mathbb{N} \times \Lambda$ the measurability requirements for the restriction of $\Gamma^{k, z}$ to $\Omega$ follow from Lemma 124 . Fix $(\vec{x}, \vec{\gamma}) \in \mathcal{Z}$. Recalling (241), if $\Gamma^{k, z}(\vec{x}, \vec{\gamma}) \neq \varnothing$, then $\operatorname{diam}_{\mathbf{X}}\left(\Gamma^{k, z}(\vec{x}, \vec{\gamma})\right) \leqslant \operatorname{diam}_{\mathbf{X}}\left(z+x_{k}+B_{\mathbf{X}}\right) \leqslant 2$. Note also that by (241) if $\Gamma^{k, z}(\vec{x}, \vec{\gamma}) \neq \varnothing$, then

$$
\Gamma^{k, z}(\vec{x}, \vec{\gamma})=\left(z+x_{k}+B_{\mathbf{X}}\right) \backslash \bigcup_{j=1}^{k-1} \bigcup_{w \in \Lambda} \Gamma^{j, w}(\vec{x}, \vec{\gamma}) .
$$

Hence $\Gamma^{k, z}(\vec{x}, \vec{\gamma}) \cap \Gamma^{j, w}(\vec{x}, \vec{\gamma})=\varnothing$ for every distinct $j, k \in \mathbb{N}$ and for every $w, z \in \Lambda$. We claim that also

$$
\Gamma^{k, z}(\vec{x}, \vec{\gamma}) \cap \Gamma^{k, w}(\vec{x}, \vec{\gamma})=\varnothing
$$

for every $k \in \mathbb{N}$ and every distinct $w, z \in \Lambda$. Indeed, it suffices to check this under the assumption that $\chi(w)=\chi(z)=\gamma_{k}$, since otherwise $\varnothing \in\left\{\Gamma^{k, z}(\vec{x}, \vec{\gamma}), \Gamma^{k, w}(\vec{x}, \vec{\gamma})\right\}$. So, suppose that

$$
\chi(w)=\chi(z)=\gamma_{k} \quad \text { yet } \quad \Gamma^{k, z}(\vec{x}, \vec{\gamma}) \cap \Gamma^{k, w}(\vec{x}, \vec{\gamma}) \neq \varnothing .
$$

By (241), this implies that there are $u, v \in B_{\mathbf{X}}$ such that $w+x_{k}+u=z+x_{k}+v$. Hence, for every $\alpha, \beta \in \mathcal{V}$,

$$
\|(w+\alpha)-(z+\beta)\|_{\mathbf{X}}=\|\alpha-\beta+v-u\|_{\mathbf{X}} \leqslant\|\alpha\|_{\mathbf{X}}+\|\beta\|_{\mathbf{X}}+\|u\|_{\mathbf{X}}+\|v\|_{\mathbf{X}} \leqslant 3+3+1+1<10,
$$

where we used the fact that $\mathcal{V} \subseteq 3 B_{\mathbf{x}}$. Since $w$ and $z$ are distinct and $\chi(w)=\chi(z)$, this is in contradiction to (240). We have thus shown that the sets $\left\{\Gamma^{k, z}(\vec{x}, \vec{\gamma})\right\}_{(k, z) \in \mathbb{N} \times \Lambda}$ are pairwise disjoint.

Note that by the definition [241], for every $(\vec{x}, \vec{\gamma}) \in \mathscr{L}$ we have

$$
\begin{equation*}
\bigcup_{k=1}^{\infty} \bigcup_{z \in \Lambda} \Gamma^{k, z}(\vec{x}, \vec{\gamma})=\bigcup_{\substack{(k, z) \in \mathbb{N} \times \Lambda \\ \chi(z)=\gamma_{k}}}\left(z+x_{k}+B_{\mathbf{X}}\right) \tag{244}
\end{equation*}
$$

Indeed, it is immediate from (241) that the left hand side of 244) is contained in the right hand side of (244). If $u$ belongs to the right hand side of (244), then let $k$ be the minimum natural number for which there is $z \in \Lambda$ with $\chi(z)=\gamma_{k}$ and $u \in z+x_{k}+B_{\mathbf{X}}$. So, for all $j \in\{1, \ldots, k-1\}$ and $w \in \Lambda$ with $\chi(w)=\gamma_{j}$ we have $u \notin w+x_{j}+B_{\mathbf{X}}$, and hence by (241) we have $v \in \Gamma^{k, z}(\vec{x}, \vec{\gamma})$, as required. By Lemma 125 if $(\vec{x}, \vec{\gamma}) \in \Omega$, then $\left\{z+x_{k}:(k, z) \in \mathbb{N} \times \Lambda \wedge \chi(z)=\gamma_{k}\right\}$ is dense in $\mathbb{R}^{n}$, and therefore the right hand side of 244 is equal to $\mathbb{R}^{n}$. Thus $\mathcal{P}$ takes values in partitions of $\mathbb{R}^{n}$.

Definition 127 introduces convenient notation that will be used several times in what follows.
Definition 127. If $\mathcal{M} \subseteq \mathbb{R}^{n}$ is Lebesgue measurable and $(k, z) \in \mathbb{N} \times \Lambda$, then define $H_{\mathcal{M}}^{k, z} \subseteq \Omega$ by

$$
\begin{equation*}
\mathrm{H}_{\mathcal{M}}^{k, z} \stackrel{\operatorname{def}}{=}\left\{(\vec{x}, \vec{\gamma}) \in \Omega: \chi(z)=\gamma_{k} \wedge z+x_{k} \in \mathcal{M}\right\} . \tag{245}
\end{equation*}
$$

IfS, $\mathcal{T} \subseteq \mathbb{R}^{n}$ are Lebesgue measurable and $(k, z) \in \mathbb{N} \times \Lambda$, then define $\mathbb{K}_{\delta, \mathcal{T}}^{k, z} \subseteq \Omega$ by

$$
\begin{equation*}
\mathrm{K}_{s, \mathcal{T}}^{k, z} \stackrel{\operatorname{def}}{=} \mathrm{H}_{\mathcal{S}}^{k, z} \backslash \bigcup_{j=1}^{k-1} \bigcup_{w \in \Lambda} \mathrm{H}_{\mathcal{T}}^{j, w} . \tag{246}
\end{equation*}
$$

The meaning of the set in (246) is that it consists of all of those $(\vec{x}, \vec{\gamma}) \in \Omega$ such that the $k$ 'th coordinate of $\vec{\gamma} \in\left\{1, \ldots, 5^{2 n}\right\}^{\mathbb{N}}$ is the color of the lattice point $z \in \Lambda$, the $k^{\prime}$ th coordinate of $\vec{x} \in \mathcal{V}^{\mathbb{N}}$ satisfies $x_{k} \in \mathcal{S}-z$, and for no $j \in\{1, \ldots, k-1\}$ and no lattice point $w \in \Lambda$ do the same assertions hold with $\mathcal{S}$ replaced by $\mathcal{T}$.
Lemma 128. Suppose that $\mathcal{S}, \mathcal{T} \subseteq \mathbb{R}^{n}$ are Lebesgue measurable sets of positive volume such that $\mathcal{S} \subseteq \mathcal{T}$. Suppose also that $\operatorname{diam}_{\mathbf{X}}(\mathcal{T}) \leqslant 4$. Then the sets

$$
\left\{\mathrm{K}_{s, J}^{k, z}\right\}_{(k, z) \in \mathbb{N} \times \Lambda}
$$

are pairwise disjoint and

$$
\begin{equation*}
\operatorname{Prob}\left[\bigcup_{k=1}^{\infty} \bigcup_{z \in \Lambda} \mathrm{~K}_{\mathcal{S}, \mathcal{T}}^{k, z}\right]=\frac{\operatorname{vol}_{n}(\mathcal{S})}{\operatorname{vol}_{n}(\mathcal{T})} \tag{247}
\end{equation*}
$$

Proof. The definition of the product measure Prob implies that for any Lebesgue measurable $\mathcal{M} \subseteq \mathbb{R}^{n}$,

$$
\begin{equation*}
\forall(j, w) \in \mathbb{N} \times \Lambda, \quad \operatorname{Prob}\left[\mathrm{H}_{\mathcal{M}}^{j, w}\right]=\mu(\mathcal{M}-w) v(\chi(w))=\frac{\operatorname{vol}_{n}(\mathcal{V} \cap(\mathcal{M}-w))}{5^{2 n_{\operatorname{vol}_{n}}(\mathcal{V})}}=\frac{\operatorname{vol}_{n}((\mathcal{V}+w) \cap \mathcal{M})}{5^{2 n_{\operatorname{vol}_{n}}(\mathcal{V})}} \tag{248}
\end{equation*}
$$

We claim if $\operatorname{diam}_{\mathbf{X}}(\mathcal{M}) \leqslant 4$, then $\left\{\mathrm{H}_{\mathcal{M}}^{j, w}\right\}_{w \in \Lambda}$ are pairwise disjoint for every fixed $j \in \mathbb{N}$. Indeed, otherwise

$$
\exists(\vec{x}, \vec{\gamma}) \in \mathrm{H}_{\mathcal{M}}^{j, w} \cap \mathrm{H}_{\mathcal{M}}^{j, z}
$$

for some distinct lattice points $w, z \in \Lambda$. Then, $\chi(w)=\gamma_{j}=\chi(z)$ and $w+x_{j}, z+x_{j} \in \mathcal{M}$. Hence,

$$
\|w-z\|_{\mathbf{X}}=\left\|\left(w+x_{j}\right)-\left(z+x_{j}\right)\right\|_{\mathbf{X}} \leqslant \operatorname{diam}_{\mathbf{X}}(\mathcal{M}) \leqslant 4 .
$$

Since $\mathcal{V} \subseteq 3 B_{\mathbf{X}}$, it follows that for every $\alpha, \beta \in \mathcal{V}$ we have

$$
\|(w+\alpha)-(z+\beta)\|_{\mathbf{X}} \leqslant\|w-z\|_{\mathbf{X}}+\|\alpha\|_{\mathbf{X}}+\|\beta\|_{\mathbf{X}} \leqslant 4+3+3=10
$$

which, by virtue of (240), contradicts the fact that $w \neq z$ and $\chi(w)=\chi(z)$.
Since $\left\{\mathrm{H}_{\mathcal{M}}^{j, w}\right\}_{w \in \Lambda}$ are pairwise disjoint and $\{w+\mathcal{V}\}_{w \in \Lambda} \operatorname{cover} \mathbb{R}^{n}$ and have pairwise disjoint interiors,

$$
\begin{equation*}
\operatorname{Prob}\left[\bigcup_{w \in \Lambda} \mathrm{H}_{\mathcal{M}}^{j, w}\right]=\sum_{w \in \Lambda} \operatorname{Prob}\left[\mathrm{H}_{\mathcal{M}}^{j, w}\right] \stackrel{\boxed{248}}{-} \frac{1}{5^{2 n} \operatorname{vol}_{n}(\mathcal{V})} \sum_{w \in \Lambda} \operatorname{vol}_{n}((\mathcal{V}+w) \cap \mathcal{M})=\frac{\operatorname{vol}_{n}(\mathcal{M})}{5^{2 n} \operatorname{vol}_{n}(\mathcal{V})} \tag{249}
\end{equation*}
$$

As $\mathcal{S} \subseteq \mathcal{T}$, we have $\operatorname{diam}_{\mathbf{X}}(\mathcal{S}) \leqslant \operatorname{diam}_{\mathbf{X}}(\mathcal{T}) \leqslant 4$. So, $\left\{\mathrm{H}_{\mathcal{S}}^{k, z_{\}}}\right\}_{z \in \Lambda}$ are pairwise disjoint for every $k \in \mathbb{N}$ by the case $\mathcal{M}=\mathcal{S}$ of the above reasoning. Recalling (246), this implies that for every $k \in \mathbb{N}$ and distinct $w, z \in \Lambda$,

$$
\mathrm{K}_{\mathcal{S}, \mathcal{T}}^{k, w} \cap \mathrm{~K}_{\mathcal{S}, \mathcal{T}}^{k, z}=\varnothing
$$

To establish that $\left\{\mathrm{K}_{\mathcal{S}, \mathcal{T}}^{k, z}\right\}_{(k, z) \in \mathbb{N} \times \Lambda}$ are pairwise disjoint it therefore remains to check that

$$
\mathrm{K}_{\mathcal{S}, \mathcal{T}}^{k, z} \cap \mathrm{~K}_{\mathcal{S}, \mathcal{T}}^{j, w}=\varnothing
$$

for every $j, k \in \mathbb{N}$ with $j<k$ and any $w, z \in \Lambda$. This is so because if $(\vec{x}, \vec{\gamma}) \in \mathrm{K}_{\mathcal{S}, \mathcal{T}}^{k, z}$, then $(\vec{x}, \vec{\gamma}) \notin \mathrm{H}_{\mathcal{T}}^{j, w}$ by (246). Therefore either $\chi(w) \neq \gamma_{j}$ or $w+x_{j} \notin \mathcal{T} \supseteq \mathcal{S}$. Consequently,

$$
(\vec{x}, \vec{\gamma}) \notin \mathrm{H}_{\mathcal{S}}^{j, w} \supseteq \mathrm{~K}_{\mathcal{S}, \mathcal{T}}^{j, w}
$$

This concludes the verification of the disjointness of $\left\{\mathrm{K}_{\mathcal{S}, \mathcal{T}}^{k, z}\right\}_{(k, z) \in \mathbb{N} \times \Lambda}$.
Since for every $k \in \mathbb{N}$ and $z \in \Lambda$, the membership of $(\vec{x}, \vec{\gamma}) \in\left\{1+, \ldots, 5^{2 n}\right\}^{\mathbb{N}} \times \mathcal{V}^{\mathbb{N}}$ in $H_{\mathcal{S}}^{k, z}$ and $H_{\mathcal{T}}^{k, z}$ depends only on the $k$ 'th coordinates of $\vec{x}$ and $\vec{\gamma}$, it follows from the independence of the coordinates that

$$
\begin{align*}
& \operatorname{Prob}\left[\mathrm{K}_{\mathcal{S}, \mathcal{T}}^{k, z}\right] \stackrel{[246}{=} \operatorname{Prob}\left[\mathrm{H}_{\mathcal{S}}^{k, z} \cap\left(\bigcap_{j=1}^{k-1}\left(\Omega \backslash \bigcup_{w \in \Lambda} \mathrm{H}_{\mathcal{T}}^{j, w}\right)\right)\right] \\
&=\operatorname{Prob}\left[\mathrm{H}_{\mathcal{S}}^{k, z}\right] \prod_{j=1}^{k-1}\left(1-\operatorname{Prob}\left[\bigcup_{w \in \Lambda} \mathrm{H}_{\mathcal{T}}^{j, w}\right]\right) \stackrel{(248}{=} \wedge \frac{\sqrt[249]{ }}{\operatorname{vol}_{n}((\mathcal{V}+z) \cap \mathcal{S})}  \tag{250}\\
& 5^{2 n} \operatorname{vol}_{n}(\mathcal{V})\left(1-\frac{\operatorname{vol}_{n}(\mathcal{T})}{\left.5^{2 n^{2} \operatorname{vol}_{n}(\mathcal{V})}\right)^{k-1}}\right.
\end{align*}
$$

Hence, since we already checked that $\left\{\mathrm{K}_{\mathcal{S}, \mathcal{T}}^{k, z}\right\}_{(k, z) \in \mathbb{N} \times \Lambda}$ are pairwise disjoint,

$$
\begin{aligned}
\operatorname{Prob}\left[\bigcup_{k=1}^{\infty} \bigcup_{z \in \Lambda} \mathrm{~K}_{\mathcal{S}, \mathcal{T}}^{k, z}\right] & =\sum_{k=1}^{\infty} \sum_{z \in \Lambda} \operatorname{Prob}\left[\mathrm{~K}_{\mathcal{S}, \mathcal{T}}^{k, z}\right] \\
& \stackrel{250}{-} \frac{1}{5^{2 n} \operatorname{vol}_{n}(\mathcal{V})}\left(\sum_{z \in \Lambda} \operatorname{vol}_{n}((\mathcal{V}+z) \cap \mathcal{S})\right) \sum_{k=1}^{\infty}\left(1-\frac{\operatorname{vol}_{n}(\mathcal{T})}{5^{2 n} \operatorname{vol}_{n}(\mathcal{V})}\right)^{k-1}=\frac{\operatorname{vol}_{n}(\mathcal{S})}{\operatorname{vol}_{n}(\mathcal{T})},
\end{aligned}
$$

where in the final step we used once more the fact that the sets $\{w+\mathcal{V}\}_{w \in \Lambda}$ cover $\mathbb{R}^{n}$ and have pairwise disjoint interiors. This completes the verification of the desired identity (247).

The following lemma is a computation of the probability of the "padding event" corresponding to the random partition $\mathcal{P}$, as a consequence of Lemma 128. In MN07] a similar argument was carried out for general finite metric spaces, but it relied on a different random partition in which the radius of the balls is also a random variable (namely, the partition of [CKR05]). This subtlety is circumvented here by using properties of normed spaces that are not available in the full generality of [MN07].

Lemma 129. Let $\mathcal{P}$ be the random partition of Lemma 126. For every $\rho \in(0,1)$ and $u \in \mathbb{R}^{n}$ we have

$$
\begin{equation*}
\operatorname{Prob}\left[u+\rho B_{\mathbf{X}} \subseteq \mathcal{P}(u)\right]=\left(\frac{1-\rho}{1+\rho}\right)^{n} . \tag{251}
\end{equation*}
$$

Proof. For every $k \in \mathbb{N}, z \in \Lambda$ and $r \in(0, \infty)$ define $\mathcal{E}_{u, r}^{k, z} \mathcal{F}_{u, r}^{k, z} \subseteq \Omega$ by

$$
\begin{equation*}
\mathcal{E}_{u, r}^{k, z} \stackrel{\text { def }}{=} \mathrm{H}_{u+r B_{\mathrm{X}}}^{k, z} \quad \text { and } \quad \mathcal{F}_{u, r}^{k, z} \stackrel{\text { def }}{=} \mathrm{K}_{u+(1-r) B_{\mathrm{X}}, u+(1+r) B_{\mathrm{X}}^{\prime}}^{k, z} \tag{252}
\end{equation*}
$$

i.e., we are using here the notations of Definition 127 for the sets $\mathcal{M}=u+r B_{\mathbf{X}}, \mathcal{S}=u+(1-r) B_{\mathbf{X}}$ and $\mathcal{T}=u+(1+r) B_{\mathbf{X}}$. We claim that

$$
\begin{equation*}
\forall(k, z) \in \mathbb{N} \times \Lambda, \quad\left\{(\vec{x}, \vec{\gamma}) \in \Omega: \Gamma^{k, z}(\vec{x}, \vec{\gamma}) \supseteq u+\rho B_{\mathbf{x}}\right\}=\mathcal{F}_{u, \rho}^{k, z} . \tag{253}
\end{equation*}
$$

Note that, since $u+(1-\rho) B_{\mathbf{X}} \subseteq u+(1+\rho) B_{\mathbf{X}}$ and $\operatorname{diam}_{\mathbf{X}}\left(u+(1+\rho) B_{\mathbf{X}}\right)=2(1+\rho) \leqslant 4$, once (253) is proven we could apply Lemma 128 to deduce the desired identity (251) as follows.

$$
\begin{aligned}
& \operatorname{Prob}\left[u+\rho B_{\mathbf{X}} \subseteq \mathcal{P}(u)\right] \stackrel{[241]}{-} \operatorname{Prob}\left[\left\{(\vec{x}, \vec{\gamma}) \in \Omega: \exists(k, z) \in \mathbb{N} \times \Lambda, \quad \Gamma^{k, z}(\vec{x}, \vec{\gamma}) \supseteq u+\rho B_{\mathbf{X}}\right\}\right] \\
& \stackrel{[253}{=} \operatorname{Prob}\left[\bigcup_{k=1}^{\infty} \bigcup_{z \in \Lambda} \mathcal{F}_{u, \rho}^{k, z}\right] \\
& \frac{247]}{} \stackrel{[252]}{ } \frac{\operatorname{vol}_{n}\left(u+(1-\rho) B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(u+(1+\rho) B_{\mathbf{X}}\right)}=\left(\frac{1-\rho}{1+\rho}\right)^{n} .
\end{aligned}
$$

To establish (253), suppose first that $(\vec{x}, \vec{\gamma}) \in \mathcal{F}_{u, \rho}^{k, z}$. By the definition of $\mathcal{F}_{u, \rho}^{k, z}$ we therefore know that

$$
\forall(j, w) \in\{1, \ldots, k-1\} \times \Lambda, \quad(\vec{x}, \vec{\gamma}) \in \mathcal{E}_{u, 1-\rho}^{k, z} \quad \text { yet } \quad(\vec{x}, \vec{\gamma}) \notin \mathcal{E}_{u, 1+\rho}^{j, w} .
$$

Hence, by the definition of $\mathcal{E}_{u, 1-\rho}^{j, w}$ we know that $\chi(z)=\gamma_{k}$ and $z+x_{k} \in u+(1-\rho) B_{\mathbf{X}}$, which (using the triangle inequality), implies that $z+x_{k}+B_{\mathbf{X}} \supseteq u+\rho B_{\mathbf{X}}$. At the same time, if $j \in\{1, \ldots, k-1\}$ and $w \in \Lambda$, then by the definition of $\mathcal{E}_{u, 1+\rho}^{j, w}$, the fact that $(\vec{x}, \vec{\gamma}) \notin \mathcal{E}_{u, 1+\rho}^{j, w}$ means that if $\chi(w)=\gamma_{j}$ then necessarily $\left\|w+x_{j}-u\right\|_{\mathbf{X}}>1+\rho$, which (using the triangle inequality) implies that $\left(w+x_{j}+B_{\mathbf{X}}\right) \cap\left(u+\rho B_{\mathbf{X}}\right)=\varnothing$. Hence, the ball $u+\rho B_{\mathbf{X}}$ does not intersect the union of the balls $\left\{w+x_{j}+B_{\mathbf{X}}:(j, w) \in\{1, \ldots, k-1\} \times \Lambda \wedge \chi(w)=\gamma_{j}\right\}$. Since $\chi(z)=\gamma_{k}$, due to (241), this implies that

$$
\Gamma^{k, z}(\vec{x}, \vec{\gamma}) \cap\left(u+\rho B_{\mathbf{X}}\right)=\left(z+x_{k}+B_{\mathbf{X}}\right) \cap\left(u+\rho B_{\mathbf{X}}\right)=u+\rho B_{\mathbf{X}},
$$

i.e., ( $\vec{x}, \vec{\gamma}$ ) belongs to the left hand side of (253).

To establish the reverse inclusion, suppose that $\Gamma^{k, z}(\vec{x}, \vec{\gamma}) \supseteq u+\rho B_{\mathbf{X}}$. The definition (241) implies in particular that $\Gamma^{k, z}(\vec{x}, \vec{\gamma}) \subseteq z+x_{k}+B_{\mathbf{X}}$ and that for $\Gamma^{k, z}(\vec{x}, \vec{\gamma})$ to be nonempty we must have $\chi(z)=\gamma_{k}$. So, we know that $\chi(z)=\gamma_{k}$ and $z+x_{k}+B_{\mathbf{X}} \supseteq u+\rho B_{\mathbf{X}}$. Assuming first that $z+x_{k} \neq u$, consider the vector

$$
\nu=u+\frac{\rho}{\left\|u-z-x_{k}\right\|_{\mathbf{X}}}\left(u-z-x_{k}\right) .
$$

Then, $v \in u+\rho B_{\mathbf{X}}$ and hence also $v \in z+x_{k}+B_{\mathbf{X}}$, i.e., $1 \geqslant\left\|v-z-x_{k}\right\|_{\mathbf{X}}=\left\|u-z-x_{k}\right\|_{\mathbf{X}}+\rho$. This shows that $\left\|z+x_{k}-u\right\|_{\mathbf{X}} \leqslant 1-\rho$, i.e., $z+x_{k} \in u+(1-\rho) B_{\mathbf{X}}$. We obtained this conclusion under the assumption that $z+x_{k} \neq u$, but it of course holds trivially also when $z+x_{k}=u$. We have thus shown that $(\vec{x}, \vec{\gamma}) \in \mathcal{E}_{u, 1-\rho}^{k, z}$.

By the definition of $\mathcal{F}_{u, p}^{k, z}$, it remains to check that

$$
\begin{equation*}
\forall(j, w) \in\{1, \ldots, k-1\} \times \Lambda, \quad(\vec{x}, \vec{\gamma}) \notin \mathcal{E}_{u, 1+\rho}^{j, w} . \tag{254}
\end{equation*}
$$

Indeed, if (254) does not hold, then let $j_{\text {min }}$ be the minimum $j \in\{1, \ldots, k-1\}$ for which $(\vec{x}, \vec{\gamma}) \in \mathcal{E}_{u, 1+\rho}^{j, w}$ for some $w \in \Lambda$. Hence, $\chi(w)=\gamma_{j_{\text {min }}}$ and $w+x_{j_{\text {min }}} \in u+(1+\rho) B_{\mathbf{X}}$. If $w+x_{j_{\text {min }}} \neq u$, then the vector

$$
u+\frac{\rho}{\left\|w+x_{j_{\min }}-u\right\|_{\mathbf{X}}}\left(w+x_{j_{\min }}-u\right)
$$

is at $\mathbf{X}$-distance $\rho$ from $u$ and also at $\mathbf{X}$-distance $\left|\rho-\left\|w+x_{j_{\text {min }}}-u\right\|_{\mathbf{X}}\right| \leqslant 1$ from $w+x_{j_{\text {min }}}$, where we used the fact that $\left\|w+x_{j_{\text {min }}}-u\right\|_{\mathbf{X}} \leqslant 1+\rho$. This shows that $\left(w+x_{j_{\text {min }}}+B_{\mathbf{X}}\right) \cap\left(u+\rho B_{\mathbf{X}}\right) \neq \varnothing$ under the assumption $w+x_{j_{\min }} \neq u$, and this assertion trivially holds also if $w+x_{j_{\text {min }}}=u$. The minimality of $j_{\min }$ implies that for every $j \in\left\{1, \ldots, j_{\min }-1\right\}$ and every $w^{\prime} \in \Lambda$ with $\chi\left(w^{\prime}\right)=\gamma_{j}$ we have $w^{\prime}+x_{j} \notin u+(1+\rho) B_{\mathrm{X}}$, i.e., $\left\|w^{\prime}+x_{j}-u\right\|_{\mathbf{X}}>1+\rho$. Hence (by the triangle inequality) we have $\left(w^{\prime}+x_{j}+B_{\mathbf{X}}\right) \cap\left(u+\rho B_{\mathbf{X}}\right)=\varnothing$. The definition of $\Gamma^{j_{\text {min }}, w}(\vec{x}, \vec{\gamma})$ now shows that $\left(u+\rho B_{\mathbf{X}}\right) \cap \Gamma^{j_{\text {min }}, w}(\vec{x}, \vec{\gamma}) \neq \varnothing$, and since by Lemma 126 we know that $\Gamma^{j_{\min }, w}(\vec{x}, \vec{\gamma})$ and $\Gamma^{k, z}(\vec{x}, \vec{\gamma})$ are disjoint (as $j_{\text {min }}<k$ ), this contradicts the premise $\Gamma^{k, z}(\vec{x}, \vec{\gamma}) \supseteq u+\rho B_{\mathbf{X}}$.

The probability of the "separation event" corresponding to the random partition $\mathcal{P}$ is estimated in the following lemma by using Lemma 128, together with input from Brunn-Minkowski theory.
Lemma 130. Let $\mathcal{P}$ be the random partition of Lemma 126 . For every $u, v \in \mathbb{R}^{n}$ we have

$$
\begin{equation*}
\operatorname{Prob}[\mathcal{P}(u) \neq \mathcal{P}(\nu)]=\min \left\{1, \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{(u-\nu)^{\perp}}\left(B_{\mathbf{X}}\right)\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}\|u-v\|_{\ell_{2}^{n}}\right\} . \tag{255}
\end{equation*}
$$

More precisely, if we denote $\psi(0)=0$ and

$$
\begin{equation*}
\forall w \in \mathbb{R}^{n} \backslash\{0\}, \quad \psi(w) \stackrel{\text { def }}{=} \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{w^{\perp}}\left(B_{\mathbf{X}}\right)\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}\|w\|_{\ell_{2}^{n}}=\frac{\|w\|_{\Pi^{*} \mathbf{X}}}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}, \tag{256}
\end{equation*}
$$

then for every $u, v \in \mathbb{R}^{n}$ we have

$$
\begin{equation*}
\frac{2 e^{\psi(u-v)}-2}{2 e^{\psi(u-v)}-1} \leqslant \operatorname{Prob}[\mathcal{P}(u) \neq \mathcal{P}(\nu)] \leqslant \frac{2 \psi(u-v)}{1+\psi(u-v)} . \tag{257}
\end{equation*}
$$

In particular, (257) implies the following more precise version of (255).

$$
\frac{2 e-2}{2 e-1} \min \{1, \psi(u-v)\} \leqslant \operatorname{Prob}[\mathcal{P}(u) \neq \mathcal{P}(\nu)] \leqslant 2 \min \{1, \psi(u-v)\} .
$$

Moreover, (257) shows that $\mathbf{P r o b}[\mathcal{P}(u) \neq \mathcal{P}(\nu)]=2 \psi(u-v)+O\left(\psi(u-v)^{2}\right)$ as $u \rightarrow v$.
Proof. If $\|u-v\|_{\mathbf{X}}>2$, then $\operatorname{Prob}[\mathcal{P}(u) \neq \mathcal{P}(\nu)]=1$ as $\mathcal{P}$ is 2-bounded. As $\left(2 e^{\psi(u-v)}-2\right) /\left(2 e^{\psi(u-v)}-1\right)<1$, the first inequality in (257) holds. By (50) we have $\psi(u-v) \geqslant\|u-v\|_{\mathbf{X}} / 2>1$, so $2 \psi(u-v) /(\psi(u-v)+1)>1$ and hence the second inequality in (257) holds. We will therefore assume from now on that $\|u-v\|_{\mathrm{X}} \leqslant 2$.

Denote $\mathcal{J}(u, v)=\left(u+B_{\mathbf{X}}\right) \cap\left(\nu+B_{\mathbf{X}}\right)$ and $\mathcal{U}(u, v)=\left(u+B_{\mathbf{X}}\right) \cup\left(v+B_{\mathbf{X}}\right)$. We claim that

$$
\begin{equation*}
\forall(k, z) \in \mathbb{N} \times \Lambda, \quad\left\{(\vec{x}, \vec{\gamma}) \in \Omega:\{u, v\} \subseteq \Gamma^{k, z}(\vec{x}, \vec{\gamma})\right\}=\mathrm{K}_{\mathcal{J}(u, v), u(u, v)}^{k, z}, \tag{258}
\end{equation*}
$$

where we recall the notation that was introduced in Definition 127. Assuming (258) for the moment, we will next explain how to conclude the proof of Lemma 130

Note that $\mathcal{J}(u, v) \subseteq \mathcal{U}(u, v)$ and $\operatorname{diam}_{\mathbf{X}}(\mathcal{U}(u, v)) \leqslant\|u-v\|_{\mathbf{X}}+2 \operatorname{diam}_{\mathbf{X}}\left(B_{\mathbf{X}}\right) \leqslant 4$. Hence, by Lemma 128 ,

$$
\begin{aligned}
\operatorname{Prob}[\mathcal{P}(u) & =\mathcal{P}(\nu)] \stackrel{[241}{=} \operatorname{Prob}\left[\left\{(\vec{x}, \vec{\gamma}) \in \Omega: \exists(k, z) \in \mathbb{N} \times \Lambda, \quad\{u, v\} \subseteq \Gamma^{k, z}(\vec{x}, \vec{\gamma})\right\}\right] \\
& \stackrel{[258}{=} \operatorname{Prob}\left[\bigcup_{k=1}^{\infty} \bigcup_{z \in \Lambda} K_{\mathcal{J}(u, v), u(u, v)}^{k, z}\right] \frac{247}{-} \frac{\operatorname{vol}_{n}(\mathcal{J}(u, v))}{\operatorname{vol}_{n}(\mathcal{U}(u, v))}=\frac{\operatorname{vol}_{n}\left(\left(u+B_{\mathbf{X}}\right) \cap\left(v+B_{\mathbf{X}}\right)\right)}{2 \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)-\operatorname{vol}_{n}\left(\left(u+B_{\mathbf{X}}\right) \cap\left(v+B_{\mathbf{X}}\right)\right)} .
\end{aligned}
$$

Hence,

$$
\begin{equation*}
\operatorname{Prob}[\mathcal{P}(u) \neq \mathcal{P}(v)]=\frac{2-2 \frac{\operatorname{vol}_{n}\left(\left(u+B_{\mathrm{X}}\right) \cap\left(v+B_{\mathrm{X}}\right)\right.}{\operatorname{vol}_{n}\left(B_{\mathrm{X}}\right)}}{2-\frac{\operatorname{vol}_{n}\left(\left(u+B_{\mathrm{X}}\right) \cap\left(v+B_{\mathrm{X}}\right)\right)}{\operatorname{vol}_{n}\left(B_{\mathrm{X}}\right)}} . \tag{259}
\end{equation*}
$$

Now, by the work [Sch92, Corollary 1] of Schmuckenschläger we have the following general estimates.

$$
\begin{equation*}
1-\psi(u-v) \leqslant \frac{\operatorname{vol}_{n}\left(\left(u+B_{\mathbf{X}}\right) \cap\left(\nu+B_{\mathbf{X}}\right)\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} \leqslant e^{-\psi(u-v)}, \tag{260}
\end{equation*}
$$

where $\psi(\cdot)$ is defined in [256). The mapping $t \mapsto(2-2 t) /(2-t)$ is decreasing on [0, 1], so 257) is consequence of 259) and 260. The remaining assertions of Lemma 130 (in particular the asymptotic evaluation (255) of the separation probability) follow from (257) by elementary calculus. Observe that for the purpose of bounding the separation modulus of $\mathbf{X}$ from above, we need only the first inequality in 260; since it is stated in [Sch92] but not proved there, for completeness we will include its elementary proof in Section 4.1.1 below. The second inequality in (260) is used here only to show that our bounds are sharp; its proof in [Sch92] relies on a more substantial use of Brunn-Minkowski theory.

It remains to verify $(258)$. $\operatorname{Fix}(k, z) \in \mathbb{N} \times \Lambda$. Suppose first that $(\vec{x}, \vec{\gamma})$ is an element of the right hand side of (258). Recalling the definitions (245) and (246), this implies that $\chi(z)=\gamma_{k}$ and $z+x_{k} \in\left(u+B_{\mathbf{X}}\right) \cap\left(v+B_{\mathbf{X}}\right)$, while for every $j \in\{1, \ldots, k-1\}$ and $w \in \Lambda$ with $\chi(w)=\gamma_{j}$ we have $w+x_{j} \notin\left(u+B_{\mathbf{X}}\right) \cup\left(\nu+B_{\mathbf{X}}\right)$. By the triangle inequality these facts imply that $z+x_{k}+B_{\mathbf{X}} \supseteq\{u, \nu\}$ and the union of the balls

$$
\left\{w+x_{j}+B_{\mathbf{X}}:(j, w) \in\{1, \ldots, k-1\} \times \Lambda \wedge \chi(w)=\gamma_{j}\right\}
$$

contains neither of the vectors $u, v$. The definition (241) of $\Gamma^{k, z}(\vec{x}, \vec{\gamma})$ now shows that $\{u, v\} \subseteq \Gamma^{k, z}(\vec{x}, \vec{\gamma})$.
For the reverse inclusion, assume that $\{u, v\} \subseteq \Gamma^{k, z}(\vec{x}, \vec{\gamma})$. Then $\chi(z)=\gamma_{k}$ and $\{u, v\} \subseteq z+x_{k}+B_{\mathbf{X}}$ by (241, which implies that $z+x_{k} \in\left(u+B_{\mathbf{X}}\right) \cap\left(v+B_{\mathbf{X}}\right)=\mathcal{J}(u, v)$. If there were $j \in\{1, \ldots, k-1\}$ and $w \in \Lambda$ with $\chi(w)=\gamma_{j}$ such that $\left(w+x_{j}+B_{\mathbf{X}}\right) \cap\{u, \nu\} \neq \varnothing$, then when one subtracts $w+x_{j}+B_{\mathbf{X}}$ from $z+x_{k}+B_{\mathbf{X}}$ one removes at least one of the vectors $u, v$, which by (241) would mean that one of these two vectors is not an element of $\Gamma^{k, z}(\vec{x}, \vec{\gamma})$, in contradiction to our assumption. Hence for all $j \in\{1, \ldots, k-1\}$ and $w \in \Lambda$ with $\chi(w)=\gamma_{j}$ we have $u \notin w+x_{j}+B_{\mathbf{X}}$ and $v \notin w+x_{j}+B_{\mathbf{X}}$, i.e., $w+x_{j} \notin\left(u+B_{\mathbf{X}}\right) \cup\left(v+B_{\mathbf{X}}\right)=U(u, v)$. This shows that $(\vec{x}, \vec{\gamma})$ belongs to the the right hand side of (258), thus completing the proof of Lemma 130 .

Proof of Theorem 121 . By rescaling, namely considering the norm $(2 / \Delta)\|\cdot\|_{\mathbf{x}}$, it suffices to treat the case $\Delta=2$. The desired random partition will then be the partition $\mathcal{P}$ of Lemma 126 and the conclusions of Theorem 121 follow from Lemma 129 and Lemma 130 .
4.1.1. Proof of the first inequality in (260). The proof of the first inequality in (260) is a simple and elementary application of standard reasoning using Fubini's theorem. Denote

$$
\begin{equation*}
t \stackrel{\text { def }}{=}\|v-u\|_{\ell_{2}^{n}} \quad \text { and } \quad x \stackrel{\text { def }}{=} \frac{1}{t}(v-u) \in S^{n-1} . \tag{261}
\end{equation*}
$$

Then,

$$
\operatorname{vol}_{n}\left(\left(u+B_{\mathbf{X}}\right) \cap\left(v+B_{\mathbf{X}}\right)\right)=\operatorname{vol}_{n}\left(B_{\mathbf{X}} \cap\left(t x+B_{\mathbf{X}}\right)\right),
$$

The desired estimate is therefore equivalent to the following assertion.

$$
\begin{equation*}
\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) \leqslant \operatorname{vol}_{n}\left(B_{\mathbf{X}} \cap\left(t x+B_{\mathbf{X}}\right)\right)+t \cdot \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}}\left(B_{\mathbf{X}}\right)\right) . \tag{262}
\end{equation*}
$$

To prove (262), partition $B_{\mathbf{X}}$ into the following three sets.

$$
\begin{align*}
& U \stackrel{\text { def }}{=} B_{\mathbf{X}} \cap\left(t x+B_{\mathbf{X}}\right),  \tag{263}\\
& V \stackrel{\text { def }}{=}\left\{y \in B_{\mathbf{X}} \backslash\left(t x+B_{\mathbf{X}}\right): \operatorname{Proj}_{x^{\perp}}(y) \in \operatorname{Proj}_{x^{\perp}}(U)\right\},  \tag{264}\\
& W \stackrel{\text { def }}{=} B_{\mathbf{X}} \backslash(U \cup V)=\left\{y \in B_{\mathbf{X}}: \operatorname{Proj}_{x^{\perp}}(y) \notin \operatorname{Proj}_{x^{\perp}}(U)\right\} . \tag{265}
\end{align*}
$$

A schematic depiction of this partition, as well as the notation of ensuing discussion, appears in Figure 3 . We recommend examining Figure 3 while reading the following reasoning because it consists of a formal justification of a situation that is clear when one keeps the geometric picture in mind.


Figure 3. A schematic depiction of the partition of $B_{\mathbf{X}}$ into the sets $U, V, W$ (with the sets $U, W$ shaded), as well as the line segments parallel to $x$ that are used in the justification of the estimate (262).

For every $z \in \operatorname{Proj}_{x^{\perp}}\left(B_{\mathbf{X}}\right)$ let $\alpha_{z} \in \mathbb{R}$ be the smallest real number such that $z+\alpha_{z} x \in B_{\mathbf{X}}$ and let $\beta_{z} \in \mathbb{R}$ be the largest real number such that $z+\beta_{z} x \in B_{\mathbf{X}}$. Thus the intersection of the line $z+\mathbb{R} x$ with $B_{\mathbf{X}}$ is the segment $w+\left[\alpha_{z}, \beta_{z}\right] x \subseteq \mathbb{R}^{n}$. Since $\|x\|_{\ell_{2}^{n}}=1$, by Fubini's theorem we have

$$
\begin{equation*}
\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)=\int_{\operatorname{Proj}_{x^{\perp}}\left(B_{X}\right)}\left(\beta_{z}-\alpha_{z}\right) \mathrm{d} z=\int_{\operatorname{Proj}_{x^{\perp}}(U)}\left(\beta_{u}-\alpha_{u}\right) \mathrm{d} u+\int_{\operatorname{Proj}_{x^{\perp}}(W)}\left(\beta_{w}-\alpha_{w}\right) \mathrm{d} w . \tag{266}
\end{equation*}
$$

For the final step of (266), note that by (265) we have $\operatorname{Proj}_{x^{\perp}}\left(B_{\mathbf{X}}\right)=\operatorname{Proj}_{x^{\perp}}(U) \cup \operatorname{Proj}_{x^{\perp}}(W)$, and the the sets $\operatorname{Proj}_{x^{\perp}}(U), \operatorname{Proj}_{x^{\perp}}(W)$ have disjoint interiors (in the subspace $x^{\perp}$ ).

Since $U=B_{\mathbf{X}} \cap\left(t x+B_{\mathbf{X}}\right)$ is convex, for every $u$ in the interior of $\operatorname{Proj}_{x^{\perp}}(U)$ the line $u+\mathbb{R} x$ intersects $U$ in an interval, say $(u+\mathbb{R} x) \cap U=u+\left[\gamma_{u}, \delta_{u}\right] x$ with $\gamma_{u}, \delta_{u} \in \mathbb{R}$ satisfying $\gamma_{u}<\delta_{u}$ such that $u+\gamma_{u} x, u+\delta_{u} x \in \partial U$ and $u+s x \in \operatorname{int}(U)$ for every $s \in\left(\gamma_{u}, \delta_{u}\right)$. Also, $(u+\mathbb{R} x) \cap B_{\mathbf{X}}=u+\left[\alpha_{u}, \beta_{u}\right] x$ with $u+\alpha_{u} x, u+\beta_{u} x \in \partial B_{\mathbf{X}}$. Thus $\left[\gamma_{u}, \delta_{u}\right] \subseteq\left[\alpha_{u}, \beta_{u}\right]$. Since $u+\gamma_{u} x \in U \subseteq t x+B_{\mathbf{X}}$, it follows that $\gamma_{w}-t \in\left[\alpha_{w}, \beta_{w}\right]$. But $\gamma_{u} \in\left[\alpha_{u}, \beta_{u}\right]$, so $\beta_{u}-\alpha_{u} \geqslant t$ and therefore $\alpha_{u}+t, \beta_{u}-t \in\left[\alpha_{u}, \beta_{u}\right]$, or equivalently $u+\left(\alpha_{u}+t\right) x, u+\left(\beta_{u}-t\right) x \in B_{\mathbf{X}}$. Because $u+\alpha_{u} x, u+\beta_{u} x \in \partial B_{\mathbf{X}}$, we get that $u+\left(\alpha_{u}+t\right) x \in B_{\mathbf{X}} \cap\left(t x+\partial B_{\mathbf{X}}\right) \subseteq \partial U$ and $u+\beta_{u} x \in\left(\partial B_{\mathbf{X}}\right) \cap\left(t x+B_{\mathbf{X}}\right) \subseteq \partial U$. Hence $\gamma_{u}=\alpha_{u}+t$ and $\delta_{u}=\beta_{u}$, from which we conclude that

$$
\begin{equation*}
u \in \operatorname{Proj}_{x^{\perp}}(U) \Longrightarrow(u+\mathbb{R} x) \cap U=u+\left[\alpha_{u}+t, \beta_{u}\right] x, \tag{267}
\end{equation*}
$$

and therefore also

$$
\begin{equation*}
u \in \operatorname{Proj}_{x^{\perp}}(U) \Longrightarrow(u+\mathbb{R} x) \cap V^{\frac{\sqrt[264]{2}}{-}} B_{\mathbf{X}} \backslash((u+\mathbb{R} x) \cap U) \stackrel{\sqrt{267}}{-} u+\left[\alpha_{u}, \alpha_{u}+t\right] x . \tag{268}
\end{equation*}
$$

Another application of Fubini's theorem now implies that

$$
\begin{align*}
\int_{\operatorname{Proj}_{x^{\perp}}(U)}\left(\beta_{u}-\alpha_{u}\right) \mathrm{d} u & =\int_{\operatorname{Proj}_{x^{\perp}}(U)} \operatorname{vol}_{1}((u+\mathbb{R} x) \cap U) \mathrm{d} u+\int_{\operatorname{Proj}_{x^{\perp}}(U)} t \mathrm{~d} u \\
& =\operatorname{vol}_{n}(U)+t \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}}(U)\right)  \tag{269}\\
& =\operatorname{vol}_{n}(U)+t\left(\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}}\left(B_{\mathbf{X}}\right)\right)-\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}}(W)\right)\right),
\end{align*}
$$

where the first step of (269) uses (267) and (268) and for the last step of (269) recall the definition (265).
Observe next that

$$
\begin{equation*}
w \in \operatorname{Proj}_{x^{\perp}}(W) \Longrightarrow \beta_{w}-\alpha_{w} \leqslant t . \tag{270}
\end{equation*}
$$

Indeed, if $w \in \operatorname{Proj}_{x^{\perp}}(W)$ yet $\beta_{w}-\alpha_{w}>t$ then $w+\left(\beta_{w}-t\right) x$ belongs to the interval joining $w+\alpha_{w} x$ and $w+\beta_{w} x$. By the convexity of $B_{\mathbf{X}}$ we therefore have $w+\left(\beta_{w}-t\right) x \in B_{\mathbf{X}}$, or equivalently $w+\beta_{w} x \in t x+B_{\mathbf{X}}$. Recalling that $w+\beta_{w} x \in B_{\mathbf{X}}$, this means that $w+\beta_{w} x \in B_{\mathbf{X}} \cap\left(t x+B_{\mathbf{X}}\right)$. By the definition 263) of $U$, it follows that $w \in \operatorname{Proj}_{x^{\perp}}(U)$. By the definition (265) of $W$, this means that $w \notin \operatorname{Proj}_{x^{\perp}}(W)$, a contradiction.

Having established (270) we see that

$$
\begin{equation*}
\int_{\operatorname{Proj}_{x^{\perp}}(W)}\left(\beta_{w}-\alpha_{w}\right) \mathrm{d} w \stackrel{\sqrt{270}}{\leqslant} t \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}}(W)\right) . \tag{271}
\end{equation*}
$$

The desired estimate (262) now follows from a substitution of 269) and 271 into 266).
4.2. Proof of Theorem 80. For any $m \in \mathbb{N}$, because $\operatorname{evr}\left(\ell_{1}^{m}\right) \asymp \sqrt{m}$, by the second part 193) of Theorem 106 there exists $\mathcal{C} \subseteq \mathbb{R}^{m}$ with $|\mathcal{C}| \leqslant e^{\beta m}$ for some universal constant $\beta>0$ such that $\mathrm{SEP}\left(\mathcal{C}_{\left.\ell_{1}^{m}\right)} \gtrsim m\right.$ (as we are considering here $\ell_{1}^{m}$ rather than more general normed spaces, this statement is due $C \mathrm{CG}^{+} 98$ ). Fix an integer $n \geqslant 2$ and $1 \leqslant p \leqslant 2$. Let $m$ be the largest integer such that $e^{\beta m} \leqslant n$. Thus $m \asymp \log n$ and

$$
\operatorname{SEP}^{n}\left(\ell_{p}\right) \geqslant \operatorname{SEP}\left(\mathcal{C}_{\ell_{p}^{m}}\right) \geqslant \frac{\operatorname{SEP}\left(\complement_{\ell_{1}^{m}}\right)}{d_{\mathrm{BM}}\left(\ell_{1}^{m}, \ell_{p}^{m}\right)} \gtrsim \frac{m}{d_{\mathrm{BM}}\left(\ell_{1}^{m}, \ell_{p}^{m}\right)}=m^{\frac{1}{p}}=(\log n)^{\frac{1}{p}}
$$

This proves the lower bound on $\operatorname{SEP}^{n}\left(\ell_{p}\right)$ in Theorem 80 .
It remains to prove the upper bound on $\operatorname{SEP}^{n}\left(\ell_{p}\right)$ in Theorem 80, i.e., that for all $x_{1}, \ldots, x_{n} \in \ell_{p}$,

$$
\begin{equation*}
\operatorname{SEP}\left(\left\{x_{1}, \ldots, x_{n}\right\},\|\cdot\|_{\ell_{p}}\right) \lesssim \frac{(\log n)^{\frac{1}{p}}}{p-1} \tag{272}
\end{equation*}
$$

The proof of 272 will refer to the following technical probabilistic lemma.
Lemma 131. Suppose that $p \in(1, \infty)$ and let $X$ be a nonnegative random variable, defined on some probability space $(\Omega$, Prob), that satisfies the following Laplace transform identity.

$$
\begin{equation*}
\forall u \in[0, \infty), \quad \mathbb{E}\left[e^{-u \mathrm{X}^{2}}\right]=e^{-u^{\frac{p}{2}}} \tag{273}
\end{equation*}
$$

Then

$$
\begin{equation*}
\mathbb{E}[\mathrm{X}]=\frac{\Gamma\left(1-\frac{1}{p}\right)}{\sqrt{\pi}}=\frac{p}{p-1} \tag{274}
\end{equation*}
$$

Moreover, we have

$$
\begin{equation*}
\forall t \in(0, \infty), \quad \operatorname{Prob}[\mathrm{X} \leqslant t] \leqslant \exp \left(-\frac{\left(\frac{p}{2}\right)^{\frac{p}{2-p}}\left(1-\frac{p}{2}\right)}{t^{\frac{2 p}{2-p}}}\right) \tag{275}
\end{equation*}
$$

Proof. Suppose that $\alpha \in(0,1)$. Then every $x \in(0, \infty)$ satisfies

$$
\begin{equation*}
\int_{0}^{\infty} \frac{1-e^{-u x}}{u^{1+\alpha}} \mathrm{d} x=x^{\alpha} \int_{0}^{\infty} \frac{1-e^{-v}}{v^{1+\alpha}} \mathrm{d} x=\frac{\Gamma(1-\alpha)}{\alpha} x^{\alpha} \tag{276}
\end{equation*}
$$

where the first step of (276) is a straightforward change of variable and the last step of (276) follows by integration by parts. The case $\alpha=1 / 2$ of (276) implies (274) as follows.

$$
\begin{aligned}
\mathbb{E}[\mathrm{X}]=\mathbb{E}\left[\frac{1}{2 \sqrt{\pi}} \int_{0}^{\infty} \frac{1-e^{-u \mathrm{X}^{2}}}{u^{\frac{3}{2}}} \mathrm{~d} u\right]= & \frac{1}{2 \sqrt{\pi}} \int_{0}^{\infty} \frac{1-\mathbb{E}\left[e^{-u \mathrm{X}^{2}}\right]}{u^{\frac{3}{2}}} \mathrm{~d} u \\
& \frac{[273)}{-} \frac{1}{2 \sqrt{\pi}} \int_{0}^{\infty} \frac{1-e^{-u^{\frac{p}{2}}}}{u^{\frac{3}{2}}} \mathrm{~d} u=\frac{1}{p \sqrt{\pi}} \int_{0}^{\infty} \frac{1-e^{-v}}{v^{1+\frac{1}{p}}} \mathrm{~d} v \frac{[276]}{-} \frac{\Gamma\left(1-\frac{1}{p}\right)}{\sqrt{\pi}} .
\end{aligned}
$$

The small ball probability estimate (275) is a consequence of the following standard use of Markov's inequality. For every $u, t \in(0, \infty)$ we have

$$
\begin{equation*}
\operatorname{Prob}[\mathrm{X} \leqslant t]=\operatorname{Prob}\left[e^{-u \mathrm{X}^{2}} \geqslant e^{-u t^{2}}\right] \leqslant e^{u t^{2}} \mathbb{E}\left[e^{-u \mathrm{X}^{2}}\right]=e^{u t^{2}-u^{\frac{p}{2}}} \tag{277}
\end{equation*}
$$

The value of $u \in(0, \infty)$ that minimizes the right-hand side of (277) is

$$
u=u(p, t) \stackrel{\text { def }}{=}\left(\frac{p}{2 t^{2}}\right)^{\frac{2}{2-p}} .
$$

A substitution of this value of $u$ into (277) simplifies to give the desired estimate (275).
Proof of (272). Fix distinct $x_{1}, \ldots, x_{n} \in \ell_{p}$. It suffices to prove (272) when $p \in(1,2)$, since the quantity that appears in the right-hand side of (272) remains bounded as $p \rightarrow 2^{-}$, and every finite subset of $\ell_{2}$ embeds isometrically into $\ell_{p}$ for every $p \in[1,2]$ (see e.g. Woj91. Chapter III.A]). We will therefore assume in the remainder of the proof of (272) that $p \in(1,2)$.

Marcus and Pisier proved [MP84, Section 2] the following statement, relying on a structural result for $p$-stable processes; its deduction from the formulation in [MP84] appears in [LMN05, Lemma 2.1]). There is a probability space ( $\Omega, \mathbf{P r o b}$ ) and a Prob-to-Borel measurable mapping $(\omega \in \Omega) \mapsto T_{\omega} \in \mathcal{L}\left(\ell_{p}, \ell_{2}\right)$ (here $\mathcal{L}\left(\ell_{p}, \ell_{2}\right)$ is the space of bounded operators from $\ell_{p}$ to $\ell_{2}$, equipped with the strong operator topology) such that for every $\omega \in \Omega$ and $x \in \ell_{p} \backslash\{0\}$ the random variable

$$
\begin{equation*}
(\omega \in \Omega) \mapsto \frac{\left\|T_{\omega}(x)\right\|_{\ell_{2}}}{\|x\|_{\ell_{p}}} \tag{278}
\end{equation*}
$$

has the same distribution as the random variable X of Lemma 131 (in particular, its distribution is independent of the choice of $\left.x \in \ell_{p} \backslash\{0\}\right)$. Consequently,

$$
\begin{equation*}
\forall i, j \in\{1, \ldots, n\}, \quad \int_{\Omega}\left\|T_{\omega}\left(x_{i}\right)-T_{\omega}\left(x_{j}\right)\right\|_{\ell_{2}} \mathrm{dProb}(\omega)=\left\|x_{i}-x_{j}\right\|_{\ell_{p}} \mathbb{E}[\mathrm{X}] \stackrel{[274}{-} \frac{\left\|x_{i}-x_{j}\right\|_{\ell_{p}}}{p-1} . \tag{279}
\end{equation*}
$$

It also follows from the above discussion and Lemma 131]that for every $t \in(0, \infty)$ we have

$$
\begin{align*}
\operatorname{Prob}[ & \left.\bigcup_{i, j \in\{1, \ldots, n\}}\left\{\omega \in \Omega:\left\|T_{\omega}\left(x_{i}\right)-T_{\omega}\left(x_{j}\right)\right\|_{\ell_{2}} \geqslant t\left\|x_{i}-x_{j}\right\|_{\ell_{p}}\right\}\right] \\
& \leqslant \sum_{\substack{i, j \in\{1, \ldots, n\} \\
i \neq j}} \operatorname{Prob}\left[\left\{\omega \in \Omega: \frac{\left\|T_{\omega}\left(x_{i}\right)-T_{\omega}\left(x_{j}\right)\right\|_{\ell_{2}}}{\left\|x_{i}-x_{j}\right\|_{\ell_{p}}}<t\right\}\right] \stackrel{[275)}{\leqslant}\binom{n}{2} \exp \left(-\frac{\left(\frac{p}{2}\right)^{\frac{p}{2-p}}\left(1-\frac{p}{2}\right)}{t^{\frac{2 p}{2-p}}}\right) . \tag{280}
\end{align*}
$$

If we choose

$$
t=t(n, p) \stackrel{\text { def }}{=} \sqrt{\frac{p}{2}}\left(\frac{2-p}{4 \log n}\right)^{\frac{1}{p}-\frac{1}{2}},
$$

then the right hand side of (280) becomes less than $1 / 2$. In other words, this shows that there exists a measurable subset $A \subseteq \Omega$ with $\operatorname{Prob}[A] \geqslant 1 / 2$ such that for every $\omega \in A$ and $i, j \in\{1, \ldots, n\}$,

$$
\begin{equation*}
\left\|x_{i}-x_{j}\right\|_{\ell_{p}} \leqslant \sqrt{\frac{2}{p}}\left(\frac{4 \log n}{2-p}\right)^{\frac{1}{p}-\frac{1}{2}}\left\|T_{\omega}\left(x_{i}\right)-T_{\omega}\left(x_{j}\right)\right\|_{\ell_{2}} \leqslant 4(\log n)^{\frac{1}{p}-\frac{1}{2}}\left\|T_{\omega}\left(x_{i}\right)-T_{\omega}\left(x_{j}\right)\right\|_{\ell_{2}}, \tag{281}
\end{equation*}
$$

where the last step of (281) uses the elementary inequality $(2 /(2-p))^{(2-p) /(2 p)} \sqrt{2 / p} \leqslant 4$, which holds (with room to spare) for every $p \in[1,2)$.
$\left\{T_{\omega}\left(x_{1}\right), \ldots, T_{\omega}\left(x_{n}\right)\right\} \subseteq \ell_{2}$ is a subset of Hilbert space of size at most $n$, so by the Johnson-Lindenstrauss dimension reduction lemma [JL84] there is $k \in \mathbb{N}$ with $k \lesssim \log n$ such that for every $\omega \in \Omega$ there is a linear operator $Q_{\omega}: \ell_{2} \rightarrow \mathbb{R}^{k}$ such that for all $i, j \in\{1, \ldots, n\}$,

$$
\begin{equation*}
\left\|T_{\omega}\left(x_{i}\right)-T_{\omega}\left(x_{j}\right)\right\|_{\ell_{2}} \leqslant\left\|Q_{\omega} T_{\omega}\left(x_{i}\right)-Q_{\omega} T_{\omega}\left(x_{j}\right)\right\|_{\ell_{2}^{k}} \leqslant 2\left\|T_{\omega}\left(x_{i}\right)-T_{\omega}\left(x_{j}\right)\right\|_{\ell_{2}} . \tag{282}
\end{equation*}
$$

An examination of the proof in [JL84] reveals that the mapping $\omega \mapsto Q_{\omega}$ can be taken to be Prob-to-Borel measurable, but actually $Q_{\omega}$ can be chosen from a finite set of operators (see e.g. Ach03).

Fix $\Delta \in(0, \infty)$. Since by $\left[\mathrm{CCG}^{+} 98\right]$ we have $\operatorname{SEP}\left(\ell_{2}^{k}\right) \lesssim \sqrt{k}$, there exists a probability space $(\Theta, \mu)$ and a mapping $\theta \mapsto \mathcal{R}^{\theta}$ that is a random partition of $\mathbb{R}^{k}$ for which

$$
\begin{equation*}
\forall(\omega, \theta, i) \in \Omega \times \Theta \times\{1, \ldots, n\}, \quad \operatorname{diam}_{\ell_{2}^{k}}\left(\mathcal{R}^{\theta}\left(Q_{\omega} T_{\omega}\left(x_{i}\right)\right)\right) \leqslant \frac{\Delta}{4(\log n)^{\frac{1}{p}-\frac{1}{2}}}, \tag{283}
\end{equation*}
$$

and also every $\omega \in \Omega$ and $i, j \in\{1, \ldots, n\}$ satisfy

$$
\begin{align*}
\mu\left(\left\{\theta \in \Theta: \mathcal{R}^{\theta}\left(Q_{\omega} T_{\omega}\left(x_{i}\right)\right) \neq \mathcal{R}^{\theta}\left(Q_{\omega} T_{\omega}\left(x_{j}\right)\right)\right\}\right) & \lesssim \frac{\sqrt{k}}{\Delta /\left(4(\log n)^{\frac{1}{p}-\frac{1}{2}}\right)}\left\|Q_{\omega} T_{\omega}\left(x_{i}\right)-Q_{\omega} T_{\omega}\left(x_{i}\right)\right\|_{\ell_{2}^{k}}  \tag{284}\\
& \lesssim \frac{(\log n)^{\frac{1}{p}}}{\Delta}\left\|T_{\omega}\left(x_{i}\right)-T_{\omega}\left(x_{i}\right)\right\|_{\ell_{2}},
\end{align*}
$$

where the last step of (284) uses the right-hand inequality in (282) and the fact that $k \lesssim \log n$.
Recalling the set $A \subseteq \Omega$ on which (281) holds for every $i, j \in\{1, \ldots, n\}$, let $v$ be the probability measure on $A$ defined by $v[E]=\operatorname{Prob}[E] / \operatorname{Prob}[A]$ for every $\operatorname{Prob}-$ measurable $E \subseteq A$ (recall that $\operatorname{Prob}[A] \geqslant 1 / 2$ ). For every $(\omega, \theta) \in A \times \Theta$ define a partition $\left.\mathcal{P}^{( } \omega, \theta\right)$ of $\left\{x_{1}, \ldots, x_{n}\right\}$ as follows.

$$
\begin{equation*}
\forall i \in\{1, \ldots, n\}, \quad \mathcal{P}^{(\omega, \theta)}\left(x_{i}\right) \stackrel{\text { def }}{=}\left\{x \in\left\{x_{1}, \ldots, x_{n}\right\}: Q_{\omega} T_{\omega}(x) \in \mathcal{R}^{\theta}\left(Q_{\omega} T_{\omega}\left(x_{i}\right)\right)\right\} . \tag{285}
\end{equation*}
$$

Then, for every $(\omega, \theta) \in A \times \Theta$ and every $i \in\{1, \ldots, n\}$ we have

$$
\begin{align*}
& \operatorname{diam}_{\ell_{p}}\left(\mathcal{P}^{(\omega, \theta)}\left(x_{i}\right)\right)=\max _{Q_{\omega} T_{\omega}\left(x_{u}\right), Q_{\omega} T_{\omega}\left(x_{\nu}\right) \in \mathbb{R}^{\theta}\left(Q_{\omega} T_{\omega}\left(x_{i}\right)\right)}\left\|x_{u}-x_{v}\right\|_{\ell_{p}} \\
& \leqslant 4(\log n)^{\frac{1}{p}-\frac{1}{2}} \quad \max _{u, v \in\{1, \ldots, n\}} \quad\left\|T_{\omega}\left(x_{u}\right)-T_{\omega}\left(x_{v}\right)\right\|_{\ell_{2}} \\
& Q_{\omega} T_{\omega}\left(x_{u}\right), Q_{\omega} T_{\omega}\left(x_{v}\right) \in \mathcal{R}^{\theta}\left(Q_{\omega} T_{\omega}\left(x_{i}\right)\right)  \tag{286}\\
& \leqslant 4(\log n)^{\frac{1}{p}-\frac{1}{2}} \quad \max _{u, v \in\{1, \ldots, n\}} \quad\left\|Q_{\omega} T_{\omega}\left(x_{u}\right)-Q_{\omega} T_{\omega}\left(x_{\nu}\right)\right\|_{\ell_{2}^{k}} \\
& Q_{\omega} T_{\omega}\left(x_{u}\right), Q_{\omega} T_{\omega}\left(x_{v}\right) \in \mathcal{R}^{\theta}\left(Q_{\omega} T_{\omega}\left(x_{i}\right)\right) \\
& \leqslant 4(\log n)^{\frac{1}{p}-\frac{1}{2}} \operatorname{diam}_{\ell_{2}^{k}}\left(\mathcal{R}^{\theta}\left(Q_{\omega} T_{\omega}\left(x_{i}\right)\right)\right) \leqslant \Delta,
\end{align*}
$$

where the first step of (286) uses (285), the second step of (286) uses 281), the third step of 286) uses (282), and the final step of 286) uses 283. Also, every distinct $i, j \in\{1, \ldots, n\}$ satisfy

$$
\begin{align*}
v \times \mu(\{(\omega, \theta) & \left.\left.\in A \times \Theta: \mathcal{P}^{(\omega, \theta)}\left(x_{i}\right) \neq \mathcal{P}^{(\omega, \theta)}\left(x_{j}\right)\right\}\right) \\
& =\int_{A} \mu\left(\left\{\theta \in \Theta: \mathcal{R}^{\theta}\left(Q_{\omega} T_{\omega}\left(x_{i}\right)\right) \neq \mathcal{R}^{\theta}\left(Q_{\omega} T_{\omega}\left(x_{j}\right)\right)\right\}\right) \mathrm{d} v(\omega) \\
& \lesssim \frac{1}{\operatorname{Prob}[A]} \int_{A} \frac{(\log n)^{\frac{1}{p}}}{\Delta}\left\|T_{\omega}\left(x_{i}\right)-T_{\omega}\left(x_{i}\right)\right\|_{\ell_{2}} \mathrm{dProb}(\omega)  \tag{287}\\
& \leqslant \frac{2(\log n)^{\frac{1}{p}}}{\Delta} \int_{\Omega}\left\|T_{\omega}\left(x_{i}\right)-T_{\omega}\left(x_{i}\right)\right\|_{\ell_{2}} \mathrm{dProb}(\omega) \\
& \lesssim \frac{(\log n)^{\frac{1}{p}}}{p-1} \cdot \frac{\left\|x_{i}-x_{j}\right\|_{\ell_{p}}}{\Delta},
\end{align*}
$$

where the first step of (287) uses (285), the second step of (287) uses (284), the third step of (287) uses $\operatorname{Prob}[A] \geqslant \frac{1}{2}$, and the last step of (287) uses 279). By (286) and (287), the proof of (272) is complete.

## 5. Barycentric-valued Lipschitz extension

In this section, we will explain how separation profiles relate to Lipschitz extension. We cannot invoke [LN05] as a "black box" because we need a more general result and our definition of random partitions differs from that of [LN05]. But, the modifications that are required in order to apply the ideas of [LN05] in the present setting are of a secondary nature, and the main geometric content of the phenomenon that is explained below is the same as in [LN05].

In addition to making the present article self-contained, there are more advantages to including here complete proofs of Theorem 65 and Theorem 113 . Firstly, the reasoning of [LN05] was designed to deal with a more general setting (treating multiple notions of random partitions at once), and it is illuminating to present a proof for separating decompositions in isolation, which leads to simplifications. Secondly, since [LN05] appeared, alternative viewpoints have been developed that relate it to optimal transport, as carried out by Kozdoba [Koz05], Brudnyi and Brudnyi [BB07a], Ohta [Oht09], and culminating more recently with a comprehensive treatment by Ambrosio and Puglisi AP20. Here we will frame the construction using the optimal transport methodology, which has conceptual advantages that go beyond yielding a clearer restructuring of the argument. The optimal transport viewpoint had an important role in quantitative improvements that were obtained in NR17, Nao21b, as well as results that will appear in forthcoming works. As a byproduct, we will use this viewpoint to easily derive a stability statement for convex hull-valued Lipschitz extension under metric transforms.
5.1. Notational preliminaries. We will start by quickly setting notation and terminology for basic concepts in measure theory and optimal transport. Everything that we describe in this subsection is standard and is included here only in order to avoid any ambiguities in the subsequent discussions.

Given a signed measure $\mu$ on a measurable space $(\Omega, \mathcal{F})$, its Hahn-Jordan decomposition is denoted $\mu=\mu^{+}-\mu^{-}$, i.e., $\mu^{+}, \mu^{-}$are disjointly supported nonnegative measures. The total variation measure of $\mu$ is $|\mu|=\mu^{+}+\mu^{-}$. For $A \in \mathcal{F}$, the restriction of $\mu$ to $A$ is denoted $\mu L_{A}$, i.e., $\mu L_{A}(E)=\mu(A \cap E)$ for $E \in \mathcal{F}$. If ( $\Omega^{\prime}, \mathcal{F}^{\prime}$ ) is another measurable space and $f: \Omega \rightarrow \Omega^{\prime}$ is a measurable mapping, then the push-forward of $\mu$ under $f$ is denoted $f_{\#} \mu$. Thus $f_{\#} \mu(E)=\mu\left(f^{-1}(E)\right)$ for $E \in \mathcal{F}^{\prime}$, or equivalently

$$
\forall h \in L_{1}\left(f_{\#} \mu\right), \quad \int_{\Omega^{\prime}} h\left(\omega^{\prime}\right) \mathrm{d} f_{\#} \mu\left(\omega^{\prime}\right)=\int_{\Omega} h(f(\omega)) \mathrm{d} \mu(\omega) .
$$

Suppose from now on that $\left(M, d_{m}\right)$ is a Polish metric space. A signed Borel measure $\mu$ on $M$ has finite first moment if $\int_{m} d_{m}(x, y) \mathrm{d}|\mu|(y)<\infty$ for all $x \in M$. Note that this implies in particular that $|\mu|(M)<\infty$,
because if $x, x^{\prime} \in T$ are distinct points, then the mapping $(y \in M) \mapsto\left[d_{m}(x, y)+d_{m}\left(x^{\prime}, y\right)\right] / d_{m}\left(x, x^{\prime}\right)$ belongs to $L_{1}(|\mu|)$ and takes values in $[1, \infty)$ by the triangle inequality.

The set of signed Borel measures on $m$ of finite first moment is denoted $\mathrm{M}_{1}\left(m, d_{m}\right)$ or simply $\mathrm{M}_{1}(m)$ if the metric is clear from the context. The set of all nonnegative measures in $\mathrm{M}_{1}(m)$ is denoted $\mathrm{M}_{1}^{+}(m)$, the set of all $\mu \in \mathrm{M}_{1}(m)$ with total mass 0 , i.e., $\mu^{+}(m)=\mu^{-}(m)$, is denoted $\mathrm{M}_{1}^{0}(m)$, and the set of all probability measures in $\mathrm{M}_{1}(m)$ is denoted $\mathrm{P}_{1}(m)$.

Given $\mu, v \in \mathrm{M}_{1}^{+}(m)$ with $\mu(m)=v(m)$, a Borel measure $\pi$ on $m \times m$ is a coupling of $\mu$ and $v$ if

$$
\pi(E \times M)=\mu(A) \quad \text { and } \quad \pi(M \times E)=v(A)
$$

for every Borel subset $E \subseteq M$. The set of couplings of $\mu$ and $v$ is denoted $\Pi(\mu, v) \subseteq \mathrm{M}_{1}^{+}(m \times m)$. Note that $(\mu \times v) / \mu(m)=(\mu \times v) / v(m) \in \Pi(\mu, v)$, so $\Pi(\mu, v) \neq \varnothing$. The Wasserstein-1 distance between $\mu$ and $v$ that is induced by the metric $d_{m}$, denoted $\mathrm{W}_{1}^{d_{m}}(\mu, v)$ or simply $\mathrm{W}_{1}(\mu, v)$ if the metric is clear from the context, is the infimum of $\int_{m \times m} d_{m}(x, y) \mathrm{d} \pi(x, y)$ over all possible couplings $\pi \in \Pi(x, y)$. Since $\left(m, d_{m}\right)$ is Polish, the metric space $\left(\mathrm{P}_{1}(m), \mathrm{W}_{1}\right)$ is also Polish; see e.g. Bol08] or AGS08, Proposition 7.1.5]. Throughout what follows, $\mathrm{P}_{1}\left(\mathrm{~m}_{)}\right.$) will be assumed to be equipped with the metric $\mathrm{W}_{1}$. The Kantorovich-Rubinstein duality theorem (see e.g. VVil09, Theorem 5.10]) asserts that

$$
\begin{equation*}
\mathrm{W}_{1}(\mu, v)=\sup _{\substack{\psi: m \rightarrow \mathbb{R} \\\|\psi\| \text { Lip }(m)=1}}\left(\int_{m} \psi \mathrm{~d} \mu-\int_{m} \psi \mathrm{~d} v\right) . \tag{288}
\end{equation*}
$$

Note that (288) implies in particular that $\mathrm{W}_{1}(\mu+\tau, v+\tau)=\mathrm{W}_{1}(\mu, v)$ for every $\tau \in \mathrm{M}_{1}^{+}(m)$.
For $\mu \in \mathrm{M}_{1}^{0}(m)$ we have $\mu^{+}(m)=\mu^{-}(m)$, so we can define $\|\mu\|_{\mathrm{W}_{1}}(m)=\mathrm{W}_{1}\left(\mu^{+}, \mu^{-}\right){ }^{13}$ This turns $\mathrm{M}_{1}^{0}(m)$ into a normed space whose completion is called the free space over $m$ (also known as the Arens-Eells space over $m$ ), and is denoted $\mathfrak{F}(m)$; see AE56, Wea99, God15] for more on this topic, and note that while $\mathfrak{F}(m)$ is commonly defined as the closure of the finitely supported measures in $\mathrm{M}_{1}^{0}\left(m_{)}\right.$) with respect to the Wasserstein-1 norm, since the finitely supported measures are dense in $\mathrm{M}_{1}^{0}(m)$ (see e.g. Vil09, Theorem 6.18]), the definitions coincide. It follows from (288) that the dual of $\mathfrak{F}(m)$ is canonically isometric to the space of all the real-valued Lipschitz functions on $m$ that vanish at some (arbitrary but fixed) point $x_{0} \in M$, equipped with the norm $\|\cdot\|_{\text {Lip }(m)}$.

Suppose that $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ is a separable Banach space and fix $\mu \in \mathrm{M}_{\mathbf{l}}(m)$. By the Pettis measurability criterion [Pet38] (see also [BL00, Proposition 5.1]), any $f \in \operatorname{Lip}(m ; Z)$ is $|\mu|$-measurable. Moreover, we have $\|f\|_{\mathbf{Z}} \in L_{1}(|\mu|)$ because if we fix $x \in M$, then for every $y \in M$,

$$
\|f(y)\|_{\mathbf{z}} \leqslant\|f(y)-f(x)\|_{\mathbf{z}}+\|f(x)\|_{\mathbf{X}} \leqslant\|f\|_{\operatorname{Lip}(m ; \mathbf{z})} d m(y, x)+\|f(x)\|_{\mathbf{X}} \in L_{1}(|\mu|),
$$

where the last step holds by the definition of $\mathrm{M}_{1}(M)$ and the fact that it implies that $|\mu|(M)<\infty$. By Bochner's integrability criterion [Boc33] (see also [BL00, Proposition 5.2]), it follows that the Bochner integrals $\int_{m} f \mathrm{~d} \mu^{+}$and $\int_{m} f \mathrm{~d} \mu^{-}$are well-defined elements of $\mathbf{Z}$, so we can consider the vector

$$
\begin{equation*}
\Im_{f}(\mu) \stackrel{\text { def }}{=} \int_{m} f \mathrm{~d} \mu=\int_{m} f \mathrm{~d} \mu^{+}-\int_{m} f \mathrm{~d} \mu^{-} \in \mathbf{Z} . \tag{289}
\end{equation*}
$$

If $\mu \in \mathrm{M}_{1}^{0}(m)$, then $\Im_{f}(\mu)=\int_{m \times m}(f(x)-f(y)) \mathrm{d} \pi(x, y)$ for every coupling $\pi \in \Pi\left(\mu^{+}, \mu^{-}\right)$. Consequently, $\left\|\widetilde{I}_{f}(\mu)\right\|_{\mathbf{Z}} \leqslant\|f\|_{\text {Lip }(m ; \mathbf{z})} \int_{m \times m} d_{m}(x, y) \mathrm{d} \pi(x, y)$, so by taking the infimum over all $\pi \in \Pi\left(\mu^{+}, \mu^{-}\right)$we see that the norm of the linear operator $\mathfrak{I}_{f}$ from $\left(\mathrm{M}_{1}^{0}(m),\|\cdot\| \mathrm{w}_{1}\right)$ to $\mathbf{Z}$ satisfies

$$
\begin{equation*}
\left\|\mathfrak{I}_{f}\right\|_{\left(M_{1}^{0}(m),\|\cdot\| w_{1}\right) \rightarrow \mathbf{Z}} \leqslant\|f\|_{\operatorname{Lip}(m ; \mathbf{z})} . \tag{290}
\end{equation*}
$$

[^9]Since $\mathrm{M}_{1}^{0}(M)$ is dense in $\mathfrak{F}(M)$, it follows that $\mathfrak{I}_{f}$ extends uniquely to a linear operator $\mathfrak{I}_{f}: \mathfrak{F}(M) \rightarrow \mathbf{Z}$ of norm at most $\|f\|_{\text {Lip }(m ; \mathbf{z})}$. So, even though elements of $\mathfrak{F}(M)$ need not be measures, one can consider the "integral" $\Im_{f}(\phi) \in \mathbf{Z}$ of $f \in \operatorname{Lip}(M ; \mathbf{Z})$ with respect to $\phi \in \mathfrak{F}(M)$; see [GK03] for more on this topic.
5.2. Refined extension moduli. Continuing with the notation that was introduced by Matoušek Mat90], we will consider the following parameters related to Lipschitz extension. Suppose that ( $M, d_{m}$ ), ( $n, d_{\eta}$ ) are metric spaces and that $\mathcal{C} \subseteq M$. Denote by e $(M, \mathcal{C} ; \eta)$ the infimum over those $K \in[1, \infty]$ such that for every $f: \mathcal{C} \rightarrow \eta$ with $\|f\|_{\operatorname{Lip}(\mathcal{C} ; \eta)}<\infty$ there is $F: m \rightarrow n$ that extends $f$ and satisfies

$$
\|F\|_{\operatorname{Lip}(m ; \eta)} \leqslant K\|f\|_{\operatorname{Lip}(\mathcal{C} ; \eta)}
$$

The supremum of $\mathrm{e}(M, \mathcal{C} ; \eta)$ over all subsets $\mathcal{C} \subseteq M$ will be denotes $\mathrm{e}(M ; n)$. Note that when $n$ is complete, $n$-valued Lipschitz functions on $\mathcal{C}$ automatically extend to the closure of $\mathcal{C}$ while preserving the Lipschitz constant, so we may assume here that $\mathcal{C}$ is closed. The supremum of $e(M, \mathcal{C} ; \mathbf{Z})$ over all Banach spaces $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ will be denoted below by e $(M, \mathcal{C})$. Thus, the notation e $(M)$ of the Introduction coincides with the supremum of $\mathrm{e}(M, \mathcal{C})$ over all subsets $\mathcal{C} \subseteq M$.

If $\left(m, d_{m}\right)$ is a metric space, $\mathcal{C} \subseteq M$, and $\left(\mathbf{Z},\|\cdot\|_{\mathbf{z}}\right)$ is a Banach space, then it is natural to consider variants of the above definitions with the additional restrictions that the extended mapping $F$ is required to take values in either the closure of the linear span of $f(\mathcal{C})$ or the closure of the convex hull of $f(\mathcal{C})$. Namely, let $\mathrm{e}_{\text {span }}(M, \mathcal{C} ; \mathbf{Z})$ be the infimum over those $K \in[1, \infty]$ such that for every $f: \mathcal{C} \rightarrow \mathbf{Z}$ there exists

$$
F: m \rightarrow \overline{\operatorname{span}}(f(\mathcal{C}))
$$

that extends $f$ and satisfies

$$
\begin{equation*}
\|F\|_{\operatorname{Lip}(m ; \mathbf{Z})} \leqslant K\|f\|_{\operatorname{Lip}\left(e^{\prime} ; \mathbf{Z}\right)} \tag{291}
\end{equation*}
$$

Analogously, let $\mathrm{e}_{\text {conv }}(M, \mathcal{C} ; \mathbf{Z})$ be the infimum over $K \in[1, \infty]$ such that for every $f: \mathcal{C} \rightarrow \mathbf{Z}$ there exists

$$
F: M \rightarrow \overline{\operatorname{conv}}(f(\mathcal{C}))
$$

that extends $f$ and satisfies (291). We then define $\mathrm{e}_{\text {conv }}(M, \mathcal{C})$ to be the supremum of $\mathrm{e}_{\text {conv }}(M, \mathcal{C} ; \mathbf{Z})$ over all possible Banach spaces $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$. Note that while one could attempt to define $\mathrm{e}_{\text {span }}(M, \mathcal{C})$ similarly, there is no point to do so because it would result in the previously defined quantity e( $M, \mathcal{C}$ ). By considering the supremum of $\mathrm{e}_{\text {conv }}(M, \mathcal{C})$ over all subsets $\mathcal{C} \subseteq M$, one defines the quantity $\mathrm{e}_{\text {conv }}(M)$.

Remark 132. By [Lin64] one can have $\mathrm{e}(M, \mathcal{C} ; \mathbf{Z})=\mathrm{e}(M ; \mathbf{Z})=1$ yet $\mathrm{e}_{\text {span }}(M, \mathcal{C}, \mathbf{Z})=\infty$ for some metric space $\left(m, d_{m}\right)$, some $\mathcal{C} \subseteq m$ and some Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{z}}\right)$. Indeed, if $\mathbf{X}$ is a closed reflexive subspace of $\ell_{\infty}$ and $\mathbf{V} \subseteq \mathbf{X}$ is a closed uncomplemented subspace of $\mathbf{X}$, then by Lin64] (see also [BL00, Corollary 7.3]) there is no Lipschitz retraction from $\mathbf{X}$ onto $\mathbf{V}$. Equivalently, the identity mapping from $\mathbf{V}$ to $\mathbf{V}$ cannot be extended to a Lipschitz mapping from $\mathbf{X}$ to $\mathbf{V}$. Hence, since $\operatorname{span}(\mathbf{V})=\mathbf{V} \subseteq \ell_{\infty}$, we have $\mathrm{e}_{\text {span }}\left(\mathbf{X}, \mathbf{V} ; \ell_{\infty}\right)=\infty$. In contrast, $\mathrm{e}\left(\mathbf{X} ; \ell_{\infty}\right)=1$ by the nonlinear Hahn-Banach theorem (see (McS34] or e.g. BL00, Lemma 1.1]). By combining [Sob41] with the discretization method of [JL84] (see also [MM16]), one can quantify the above example by showing that for arbitrarily large $n \in \mathbb{N}$ there are Banach spaces $\left(\mathbf{X},\|\cdot\|_{\mathbf{X}}\right),\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ and a subset $\mathcal{C} \subseteq \mathbf{X}$ with $|\mathcal{C}|=n$ for which we have

$$
\begin{equation*}
\frac{\mathrm{e}_{\text {span }}(\mathbf{X}, \mathcal{C} ; \mathbf{Z})}{\mathrm{e}(\mathbf{X}, \mathcal{C} ; \mathbf{Z})} \gtrsim \sqrt{\frac{\log n}{\log \log n}} \tag{292}
\end{equation*}
$$

(In fact, in (292) one can have $\mathrm{e}(\mathbf{X}, \mathcal{C} ; \mathbf{Z})=\mathrm{e}(\mathbf{X} ; \mathbf{Z})=1$.) At present, the right hand side of 292 is the largest asymptotic dependence on $n$ that we are able to obtain for this question, and it remains an interesting open problem to determine the best possible asymptotics here.

Most, but not all, of the Lipschitz extension methods in the literature, including Kirszbraun's extension theorem [Kir34], Ball's extension theorem Bal92] and methods that rely on (variants of) partitions of unity such as in [JLS86, LN05, LS05, BB06], yield convex hull-valued extensions, i.e., they actually provide bounds on the quantity $\mathrm{e}_{\mathrm{conv}}(M, \mathcal{C} ; \mathbf{Z})$. Nevertheless, it seems likely that there is no $\varphi:[1, \infty) \rightarrow[1, \infty)$
such that $\mathrm{e}_{\text {conv }}(M) \leqslant \varphi(\mathrm{e}(M))$ for every Polish metric space $\left(M, d_{m}\right)$, though if such an estimate were available, then it would be valuable; see e.g. Remark 140 . In fact, we propose the following conjecture.

Conjecture 133. There exists a Polish metric space $\left(M, d_{m}\right)$ for which $\mathrm{e}(M)<\infty$ yet $\mathrm{e}_{\mathrm{conv}}(M)=\infty$.
Remark 134. By definition, for every metric space $\left(M, d_{m}\right)$, every Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{z}}\right)$ and every $\mathcal{C} \subseteq M$,

$$
\mathrm{e}_{\text {conv }}(M, \mathrm{C} ; \mathbf{Z}) \geqslant \mathrm{e}_{\text {span }}(M, \mathcal{C} ; \mathbf{Z}) \geqslant \mathrm{e}(M, \mathcal{C} ; \mathbf{Z})
$$

We explained in Remark 132 that the second of these inequalities can be strict (in a strong sense). However, as a complement to Conjecture 133, we state that to the best of our knowledge it is unknown whether this is so for the first of these inequalities, i.e., if it could happen that $\mathrm{e}_{\text {span }}(M, \mathcal{C} ; \mathbf{Z})<\infty$ yet $\mathrm{e}_{\text {conv }}(M, \mathcal{C} ; \mathbf{Z})=\infty$. We suspect that this is possible, but if not, then it would be interesting to investigate how one could bound $\mathrm{e}_{\text {conv }}(M, \mathcal{C} ; \mathbf{Z})$ from above by a function of $\mathrm{e}_{\text {span }}(M, \mathcal{C} ; \mathbf{Z})$. We do know that there are a metric space $\left(M, d_{m}\right)$, a Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{z}}\right)$, a subset $\mathcal{C} \subseteq M$ and a Lipschitz mapping $f: \mathcal{C} \rightarrow \mathbf{Z}$ that can be extended to a Lipschitz mapping that takes values in $\overline{\operatorname{span}}(f(\mathcal{C}))$ but cannot be extended to a Lipschitz mapping that takes values in $\overline{\operatorname{conv}}\left(f(\mathcal{C})\right.$. To see this, let $\left\{e_{j}\right\}_{j=1}^{\infty}$ be the standard basis of $\ell_{\infty}$. For $n \in \mathbb{N}$ set $m(n)=n(n-1) / 2$ and let $\mathbf{X}_{n}$ be the span of $\left\{e_{m(n)+1}, \ldots, e_{m(n+1)}\right\}$ in $\ell_{\infty}$. Thus, $\mathbf{X}_{n}$ is isometric to $\ell_{\infty}^{n}$ and $\ell_{\infty}=\left(\oplus_{n=1}^{\infty} \mathbf{X}_{n}\right)_{\infty}$. By [Sob41], there is a linear subspace $\mathbf{V}_{n}$ of $\mathbf{X}_{n}$ such that every linear projection $Q: \mathbf{X}_{n} \rightarrow \mathbf{V}_{n}$ satisfies $\|Q\|_{\mathbf{X}_{n} \rightarrow \mathbf{V}_{n}} \gtrsim \sqrt{n}$. By the method of [JL84], it follows that there exists ${ }^{14}$ $\mathcal{A}_{n} \subseteq B_{V_{n}}=V_{n} \cap B_{\ell_{\infty}}$ with $\left|\mathcal{A}_{n}\right| \leqslant n^{O(n)}$ such that $\left\|F_{n}\right\|_{\operatorname{Lip}\left(X_{n} ; V_{n}\right)} \gtrsim \sqrt{n}$ for any $F_{n}: \mathbf{X}_{n} \rightarrow \mathbf{V}_{n}$ that extends the formal identity $\operatorname{ld}_{\mathcal{A}_{n} \rightarrow \mathbf{V}_{n}}: \mathcal{A}_{n} \rightarrow \mathbf{V}_{n}$. By compactness, there exists $\delta_{n} \in(0,1)$ such that if we define

$$
\mathcal{C}_{n}=\mathcal{A}_{n} \cup\left\{\delta_{n} e_{m(n)+1}, \ldots, \delta_{n} e_{m(n+1)}\right\} \cup\{0\}
$$

then also $\left\|\Phi_{n}\right\|_{\operatorname{Lip}\left(\mathbf{x}_{n} ; \mathbf{X}_{n}\right)} \gtrsim \sqrt{n}$ for any mapping $\Phi_{n}$ from $\mathbf{X}_{n}$ to the polytope $\overline{\operatorname{conv}}\left(\mathcal{C}_{n}\right)$ that extends the formal identity $\operatorname{ld}_{\mathrm{C}_{n} \rightarrow \mathbf{x}_{n}}$. Consider the subset

$$
\mathcal{C}=\bigcup_{n=1}^{\infty} \mathcal{C}_{n} \subseteq \ell_{\infty}
$$

If $\Phi: \ell_{\infty} \rightarrow \overline{\operatorname{conv}}(\mathcal{C})$ extends $\mathrm{Id}_{\mathcal{C} \rightarrow \ell_{\infty}}$, then for each $n \in \mathbb{N}$ the mapping $\mathrm{R}_{n} \circ\left(\left.\Phi\right|_{\mathbf{X}_{n}}\right): \mathbf{X}_{n} \rightarrow \mathbf{X}_{n}$ extends $\operatorname{Id}_{\mathcal{C}_{n} \rightarrow \ell_{\infty}}$ and takes values in $\overline{\operatorname{conv}}\left(\mathcal{C}_{n}\right)$, where $\mathrm{R}_{n}: \ell_{\infty} \rightarrow \mathbf{X}_{n}$ is the canonical restriction operator. Hence,

$$
\|\Phi\|_{\operatorname{Lip}\left(\ell_{\infty} ; \mathbf{x}_{n}\right)} \geqslant\left\|\mathrm{R}_{n} \circ\left(\left.\Phi\right|_{\mathbf{x}_{n}}\right)\right\|_{\operatorname{Lip}\left(\mathbf{X}_{n} ; \mathbf{x}_{n}\right)} \gtrsim \sqrt{n}
$$

Since this holds for every $n \in \mathbb{N}$, the mapping $\Phi$ is not Lipschitz. Consequently, $\mathrm{e}_{\text {conv }}\left(\ell_{\infty}, \mathcal{C} ; \ell_{\infty}\right)=\infty$. At the same time, by construction we have $\overline{\operatorname{span}}(\mathcal{C})=\overline{\operatorname{span}}\left(\left\{e_{j}\right\}_{j=1}^{\infty}\right)=c_{0}$ (recall that $c_{0}$ commonly denotes the subspace of $\ell_{\infty}$ consisting of all those sequences that tend to 0 ). So, any 2 -Lipschitz retraction $\rho$ of $\ell_{\infty}$ onto $c_{0}$ extends $\mathrm{Id}_{\mathrm{C}_{\rightarrow} \ell_{\infty}}$ and takes values in $\overline{\operatorname{span}}(\mathcal{C})$; the existence of such a retraction $\rho$ is due to Lin64] (see also BL00, Example 1.5]). If $\mathrm{e}_{\text {span }}\left(\ell_{\infty}, \mathcal{C} ; \ell_{\infty}\right)$ were finite, then this example would answer the above question, ${ }^{15}$ but we suspect that in fact $\mathrm{e}_{\text {span }}\left(\ell_{\infty}, \mathcal{C} ; \ell_{\infty}\right)=\infty$.

Proposition 135 is a convenient characterization of the quantities e( $M, \mathcal{C}$ ) and $\mathrm{e}_{\operatorname{conv}}(\mathcal{M}, \mathcal{C})$; while it was not previously stated explicitly in this form, its proof is based on well-understood ideas.

Proposition 135. Suppose that $\left(M, d_{m}\right)$ is a metric space, $\mathcal{C}$ is a Polish subset of $M$ and $s_{0} \in \mathcal{C}$. Fix two nonnegative functions $\mathfrak{d}: M \times M \rightarrow[0, \infty)$ and $\varepsilon: \mathcal{C}: \rightarrow[0, \infty)$. Then, the following two equivalences hold.
(1) The following two statements are equivalent.

[^10]- For every Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ and every mapping $f: \mathrm{C} \rightarrow \mathbf{Z}$ that is 1 -Lipschitz with respect to the metric $d_{m}$ there exists $F: m \rightarrow \mathbf{Z}$ that satisfies the following two conditions.
$-\|F(s)-f(s)\|_{\mathbf{Z}} \leqslant \varepsilon(s)$ for every $s \in \mathcal{C}$.
- $\|F(x)-F(y)\|_{\mathbf{z}} \leqslant \mathfrak{d}(x, y)$ for every $x, y \in W$.
- There exists a family $\left\{\phi_{x}\right\}_{x \in}$ of elements of the free space $\mathfrak{F}(\mathcal{C})$ with the following properties.
$-\left\|\phi_{s}-\boldsymbol{\delta}_{s}+\boldsymbol{\delta}_{s_{0}}\right\|_{\mathcal{F}(\mathrm{C})} \leqslant \varepsilon(s)$ for every $s \in \mathcal{C}$.
- $\left\|\phi_{x}-\phi_{y}\right\|_{\mathfrak{F}(\mathrm{e})} \leqslant \mathfrak{d}(x, y)$ for every $x, y \in T$.
(2) The following two statements are equivalent.
- For every Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ and every mapping $f: \mathrm{C} \rightarrow \mathbf{Z}$ that is 1-Lipschitz with respect to the metric $d_{m}$ there exists $F: m \rightarrow \overline{\operatorname{conv}}(f(\mathcal{C}))$ that satisfies the following two conditions.
$-\|F(s)-f(s)\|_{\mathbf{Z}} \leqslant \varepsilon(s)$ for every $s \in \mathcal{C}$.
- $\|F(x)-F(y)\|_{\mathbf{z}} \leqslant \mathfrak{d}(x, y)$ for every $x, y \in T$.
- There exists a family $\left\{\mu_{x}\right\}_{x \in m}$ of probability measures in $\mathrm{P}_{1}(\mathrm{C})$ with the following properties.
$-\mathrm{W}_{1}^{d_{m}}\left(\mu_{s}, \boldsymbol{\delta}_{s}\right) \leqslant \varepsilon(s)$ for every $s \in \mathcal{C}$.
- $\mathrm{W}_{1}^{d_{m}}\left(\mu_{x}, \mu_{y}\right) \leqslant \mathfrak{d}(x, y)$ for every $x, y \in T$.

In the setting of Proposition 135, if $\varepsilon(s)=0$ for every $s \in \mathcal{C}$ and also $\mathfrak{d}=K d_{m}$ for some $K \geqslant 1$, then in AP20, Definition 2.7] a family $\left\{\phi_{x}\right\}_{x \in} \subseteq \mathfrak{F}(\mathcal{C})$ as in part (1) of Proposition 135 is called a $K$-random projection of $M$ onto $\mathcal{C}$, and in [Oht09, Definition 3.1] a family $\left\{\mu_{x}\right\}_{x \in m} \subseteq \mathrm{P}_{1}(\mathcal{C})$ as in part (2) of Proposition 135 is called a stochastic $K$-Lipschitz retraction of $M$ onto $\mathcal{C}$ while in [AP20, Definition 2.7] it is called a strong $K$-random projection of $m$ onto $\mathcal{C}$.

Proof of Proposition 135. Suppose first that $\left\{\phi_{x}\right\}_{x \in} \subseteq \mathfrak{F}(\mathcal{C})$ and $\left\{\mu_{x}\right\}_{x \in M} \subseteq \mathrm{P}_{1}(\mathcal{C})$ are as in the two parts of Proposition 135. Let $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ be a Banach space and fix a 1 -Lipschitz function $f: \mathcal{C} \rightarrow \mathbf{Z}$. Since $\mathcal{C}$ is Polish and hence separable, by replacing $\mathbf{Z}$ with the closure of the linear span of $f(\mathcal{C})$ we may assume that $\mathbf{Z}$ is separable. Recalling the notation (289) and the discussion immediately following it for the (integration) operator $\mathfrak{I}_{f}: \mathrm{M}_{1}(m) \cup \mathfrak{F}(m) \rightarrow \mathbf{Z}$, define two (linear) mappings

$$
\mathrm{Ext}_{\mathrm{e}}^{\phi} f, \mathrm{Ext}_{\mathrm{e}}{ }^{\mu} f: m \rightarrow \mathbf{Z}
$$

by setting for every $x \in m$,

$$
\begin{equation*}
\mathrm{Ext}_{\mathrm{e}}^{\phi} f(x) \stackrel{\text { def }}{=} f\left(x_{0}\right)+\Im_{f}\left(\phi_{x}\right) \quad \text { and } \quad \mathrm{Ext}_{\mathrm{e}}^{\mu} f(x) \stackrel{\text { def }}{=} \Im_{f}\left(\mu_{x}\right) \stackrel{289}{=} \int_{\mathrm{e}} f \mathrm{~d} \mu_{x} . \tag{293}
\end{equation*}
$$

Observe that since $\mu_{x}$ is a probability measure, Ext ${ }_{\mathrm{C}}^{\mu} f(x)$ belongs to the closure of the convex hull of $f(\mathcal{C})$.
For every $x, y \in M$ we have

$$
\left\|\operatorname{Ext}_{\mathcal{C}}^{\phi} f(x)-\operatorname{Ext}_{\mathfrak{C}}^{\phi} f(y)\right\|_{\mathbf{Z}}=\left\|\mathfrak{I}_{f}\left(\phi_{x}-\phi_{y}\right)\right\|_{\mathbf{Z}} \stackrel{\sqrt{290}}{\leqslant}\left\|\phi_{x}-\phi_{y}\right\|_{\mathfrak{F}(\mathcal{C})} \leqslant \mathfrak{d}(x, y),
$$

and similarly (using Kantorovich-Rubinstein duality),

$$
\left\|\operatorname{Ext}_{\mathrm{e}}^{\mu} f(x)-\operatorname{Ext}_{\mu}^{\phi} f(y)\right\|_{\mathbf{Z}} \leqslant \mathrm{W}_{1}^{d_{m}}\left(\mu_{x}, \mu_{y}\right) \leqslant \mathfrak{d}(x, y)
$$

Also, for every $s \in \mathcal{C}$ we have

$$
\left\|\operatorname{Ext}^{\mathcal{C}}{ }^{\phi} f(s)-f(s)\right\|_{\mathbf{Z}}=\left\|\Im_{f}\left(\phi_{s}-\boldsymbol{\delta}_{s}+\boldsymbol{\delta}_{x_{0}}\right)\right\|_{\mathbf{Z}} \leqslant\left\|\phi_{s}-\boldsymbol{\delta}_{s}+\boldsymbol{\delta}_{x_{0}}\right\|_{\mathfrak{F}(\mathcal{C})} \leqslant \varepsilon(s),
$$

and similarly,

$$
\left\|\operatorname{Ext}_{\mathrm{C}}^{\mu} f(s)-f(s)\right\|_{\mathbf{Z}}=\left\|\Im_{f}\left(\phi_{s}-\boldsymbol{\delta}_{s}\right)\right\|_{\mathbf{Z}} \leqslant \mathrm{W}_{1}^{d_{m}}\left(\mu_{s}, \boldsymbol{\delta}_{s}\right) \leqslant \varepsilon(s) .
$$

Conversely, define $f: \mathcal{C} \rightarrow \mathfrak{F}(\mathcal{C})$ by setting $f(s)=\boldsymbol{\delta}_{s}-\boldsymbol{\delta}_{s_{0}}$ for each $s \in \mathcal{C}$. Then $f$ is 1-Lipschitz. Fix $F: M \rightarrow \mathfrak{F}(\mathcal{C})$. Writing $F(x)=\phi_{x}$ for each $x \in M$, the assumptions of the first half of part (1) of Proposition 135 coincide with the assertions of its second half. As $\mathcal{C}$ is Polish, $\mathrm{P}_{1}(\mathcal{C})$ is closed in $\mathfrak{F}(\mathcal{C})$. Therefore,

$$
\overline{\operatorname{conv}}(f(\mathcal{C}))=\mathrm{P}_{1}(\mathcal{C})-\boldsymbol{\delta}_{s_{0}},
$$

where the closure is in $\mathfrak{F}(\mathcal{C})$. Thus, if $F(M) \subseteq \overline{\operatorname{conv}}(f(\mathcal{C}))$, then $\mu_{x} \stackrel{\text { def }}{=} F(x)+\boldsymbol{\delta}_{s_{0}} \in \mathrm{P}_{1}(\mathcal{C})$ and the assumptions of the first half of part (2) of Proposition 135 coincide with the assertions of its second half.

The proof of Proposition 135 shows that even though in the first parts of the two equivalences in Proposition 135 one assumes merely the existence of an $F$ with the desired properties, it follows that such an $F$ can in fact be chosen to depend linearly on the input $f$, per 293.

Due to Proposition 135, the following question is closely related to Conjecture 133 , though we think that it is also of independent interest.

Question 136. Characterize those Polish metric spaces ( $M, d_{m}$ ) for which there exists a Lipschitz mapping $\rho: \mathfrak{F}(M) \rightarrow \mathrm{P}_{1}(M)$ (recall that by default $\mathrm{P}_{1}(M)$ is equipped with the Wasserstein-1 metric) and $x_{0} \in M$ such that $\rho\left(\delta_{y}-\delta_{x_{0}}\right)=\delta_{y}$ for every $y \in M$.
5.3. Barycentric targets. Following [MN13], say that a metric space ( $M, d_{m}$ ) is $\mathrm{W}_{1}$-barycentric with constant $\beta>0$ if there is a mapping $\mathfrak{B}: \mathrm{P}_{1}(M) \rightarrow m$ that satisfies $\mathfrak{B}\left(\boldsymbol{\delta}_{x}\right)=x$ for every $x, \in M$, and also

$$
\forall \mu, v \in \mathrm{P}_{1}(m), \quad d_{m}(\mathfrak{B}(\mu), \mathfrak{B}(v)) \leqslant \beta \mathrm{W}_{1}^{d_{m}}(\mu, v)
$$

The infimal $\beta$ for which this holds is denoted $\beta_{1}(M)$. This notion (and variants thereof) were studied in various contexts; see e.g. (ESH99, LPS00, Gro03, Stu03, LN05, Oht09, Aus11, MN13, Nav13, Lim18, Bas18. Any normed space $\mathbf{X}$ is $W_{1}$-barycentric with constant 1 , as seen by considering $\mathfrak{B}(\mu)=\int_{X} x \mathrm{~d} \mu(x)$. Other examples of spaces that are $W_{1}$-barycentric with constant 1 include Hadamard spaces and Busemann nonpositively curved spaces [BH99], or more generally spaces with a conical geodesic bicombing [DL15].

Thanks to Proposition 135, convex hull-valued (approximate) extension theorems automatically generalize to extension theorems for mappings that take value in $W_{1}$-barycentric metric spaces.

Proposition 137. Let $\left(M, d_{m}\right)$ be a metric space and suppose that $\mathcal{C} \subseteq M$ is a Polish subset of $M$. Fix $\mathfrak{d}: M \times M \rightarrow[0, \infty)$ and $\varepsilon: \mathcal{C} \rightarrow[0, \infty)$. Assume that for every Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ and every $f: \mathcal{C} \rightarrow \mathbf{Z}$ that is 1-Lipschitz with respect to $d_{m}$ there is $F: M \rightarrow \overline{\operatorname{conv}}(f(\mathcal{C}))$ that satisfies $\|F(s)-f(s)\|_{\mathbf{z}} \leqslant \varepsilon(s)$ for every $s \in \mathcal{C}$ and $\|F(x)-F(y)\|_{\mathbf{z}} \leqslant \mathfrak{d}(x, y)$ for every $x, y \in M$. Fix $\eta: \mathcal{C} \rightarrow(1, \infty)$ and $\tau: M \times M \rightarrow(1, \infty)$, as well as $\beta>0$ and a concave nondecreasing function $\omega:[0, \infty) \rightarrow[0, \infty)$ with $\omega(0)=0$. If $\left(n, d_{n}\right)$ is a $\mathrm{W}_{1}$-barycentric metric space with constant $\beta$ and $\phi: \mathcal{C} \rightarrow \eta$ has modulus of uniform continuity $\omega$ with respect to $d_{m}$, namely $d_{n}(f(s), f(t)) \leqslant \omega\left(d_{m}(s, t)\right)$ for every $s, t \in \mathcal{C}$, then there is $\Phi: m \rightarrow \eta$ such that $d_{n}(\Phi(s), \phi(s)) \leqslant \omega(\eta(s) \varepsilon(s))$ for every $s \in \mathcal{C}$ and $d_{n}(\Phi(x), \Phi(y)) \leqslant \omega(\tau(x, y) \mathfrak{d}(x, y))$ for every $x, y \in M$.
Proof. By Proposition 135 , there is a collection of measures $\left\{\mu_{x}\right\}_{x \in} \subseteq \mathrm{P}_{1}(\mathrm{C})$ such that

$$
\forall s \in \mathcal{C}, \quad \mathrm{~W}_{1}^{d_{m}}\left(\mu_{s}, \boldsymbol{\delta}_{s}\right) \leqslant \varepsilon(s) \quad \text { and } \quad \forall x, y \in m \quad \mathrm{~W}_{1}^{d_{m}}\left(\mu_{x}, \mu_{y}\right) \leqslant \mathfrak{d}(x, y)
$$

Hence, for every $s \in \mathcal{C}$ and $x, y \in M$ there are couplings $\pi_{s} \in \Pi\left(\mu_{s}, \boldsymbol{\delta}_{s}\right)$ and $\pi_{x, y} \in \Pi\left(\mu_{x}, \mu_{y}\right)$ such that

$$
\iint_{\mathfrak{C} \times \mathcal{C}} d_{m}(u, v) \mathrm{d} \pi_{s}(u, v) \leqslant \eta(s) \varepsilon(s) \quad \text { and } \quad \iint_{\mathfrak{C} \times \mathfrak{C}} d_{m}(u, v) \mathrm{d} \pi_{x, y}(u, v) \leqslant \tau(x, y) \mathfrak{d}(x, y)
$$

Since $(\phi \times \phi)_{\#} \pi_{s} \in \Pi\left(\phi_{\#} \mu_{s}, \phi_{\#} \boldsymbol{\delta}_{s}\right)$ and $(\phi \times \phi)_{\#} \pi_{x, y} \in \Pi\left(\phi_{\#} \mu_{x}, \phi_{\#} \mu_{y}\right)$, it follows that

$$
\begin{aligned}
\mathrm{W}_{1}^{d_{n}}\left(\phi_{\#} \mu_{s}, \phi_{\#} \boldsymbol{\delta}_{s}\right) & \leqslant \iint_{n \times n} d_{n}(a, b) \mathrm{d}(\phi \times \phi)_{\#} \pi_{s}(a, b) \\
& =\iint_{n \times n} d_{n}(\phi(u), \phi(v)) \mathrm{d} \pi_{s}(u, v) \\
& \leqslant \iint_{n \times n} \omega\left(d_{n}(u, v)\right) \mathrm{d} \pi_{s}(u, v) \\
& \leqslant \omega\left(\iint_{n \times n} d_{n}(u, v) \mathrm{d} \pi_{s}(u, v)\right) \\
& \leqslant \omega(\eta(s) \varepsilon(s))
\end{aligned}
$$

where the penultimate step uses the concavity of $\omega$. For the same reason, also

$$
\mathrm{W}_{1}^{d_{n}}\left(\phi_{\#} \mu_{x}, \phi_{\#} \mu_{y}\right) \leqslant \omega(\tau(x, y) \mathfrak{d}(x, y)) .
$$

Since $\left(\eta, d_{n}\right)$ is $\beta$-barycentric there is $\mathfrak{B}: \mathrm{P}_{1}(\eta) \rightarrow \eta$ satisfying $\mathfrak{B}\left(\boldsymbol{\delta}_{z}\right)=z$ for every $z, \epsilon n$, and

$$
\forall v_{1}, v_{2} \in \mathrm{P}_{1}(\eta), \quad d_{n}\left(\mathfrak{B}\left(v_{1}\right), \mathfrak{B}\left(v_{2}\right)\right) \leqslant \beta \mathrm{W}_{1}^{d_{n}}\left(v_{1}, v_{2}\right) .
$$

Define $\Phi: m \rightarrow n$ by

$$
\forall x \in m, \quad \Phi(x) \stackrel{\text { def }}{=} \mathfrak{B}\left(\phi_{\#} \mu_{x}\right) .
$$

Then, for every $s \in \mathcal{C}$ we have

$$
d_{n}(\Phi(s), \phi(s)) \leqslant \beta \mathrm{W}_{1}^{d_{n}}\left(\phi_{\#} \mu_{s}, \phi_{\#} \boldsymbol{\delta}_{s}\right) \leqslant \omega(\eta(s) \varepsilon(s)),
$$

and for the same reason also $d_{\eta}(\Phi(x), \phi(y)) \leqslant \omega(\tau(x, y) \mathfrak{d}(x, y))$ for every $x, y \in m$.
Because (as we will soon see) all of our new Lipschitz extension theorems are in fact bounds on $\mathrm{e}_{\text {conv }}(\cdot)$, the following immediate corollary of Proposition 137 (with $\mathfrak{d}$ a multiple of $d_{m}$ and $\omega$ linear) shows that they apply to barycentric targets and not only to Banach space targets.
Corollary 138. Fix $\beta>0$. Suppose thatTh is a Polish metric space and that $h$ is a complete $\mathrm{W}_{1}$-barycentric metric space with constant $\beta$. Then, $\mathrm{e}_{\text {conv }}(m, n) \leqslant \beta \mathrm{e}_{\text {conv }}(m)$.

Another noteworthy special case of Proposition 137 is when $\omega(s)=s^{\theta}$ for some $0<\theta \leqslant 1$, i.e., in the setting of Hölder extension that we discussed in Remark 15 and Section 2.3 . Analogously to (18), we denote the convex hull-valued $\theta$-Hölder extend modulus of a metric space $\left(m, d_{m}\right)$ by

$$
\mathrm{e}_{\mathrm{conv}}^{\theta}(m)=\mathrm{e}_{\mathrm{conv}}\left(m, d_{m}^{\theta}\right) .
$$

Corollary 139. Suppose that $m$ is a Polish metric space. Then, for every $0<\theta \leqslant 1$ we have

$$
\mathrm{e}^{\theta}(m) \leqslant \mathrm{e}_{\text {conv }}^{\theta}(m) \leqslant \mathrm{e}_{\text {conv }}(m)^{\theta} .
$$

Because the upper bound on $\mathrm{e}\left(\ell_{\infty}^{n}\right)$ that we obtain in Theorem 14 is actually an upper bound on $\mathrm{e}_{\text {conv }}\left(\ell_{\infty}^{n}\right)$, Corollary (139) implies (19). More generally, Proposition 137 implies that

$$
\mathrm{e}_{\mathrm{conv}}(m, \omega \circ d m) \leqslant \sup _{d>0} \frac{\omega\left(\mathrm{e}_{\mathrm{conv}}(m) d\right)}{\omega(d)}
$$

for any concave nondecreasing function $\omega:[0, \infty) \rightarrow[0, \infty)$ with $\omega(0)=0$.
Remark 140. The question of how Lipschitz extension results imply extension results for other moduli of uniform continuity was studied in [Nao01] and treated definitively by Brudnyi and Shvartsman in BS02 using an interesting connection to the Brudny̆-Krugljak $K$-divisibility theorem [BK81] (see also [Cwi84]) from the theory of real interpolation of Banach spaces. In particular, by [BS02] we have e ${ }^{\theta}(m) \lesssim \mathrm{e}(m)^{2}$, which remains the best-known bound on $\mathrm{e}^{\theta}(m)$ in terms of $\mathrm{e}\left(m_{)}\right.$and it would be interesting to determine if it could be improved. As Corollary 139 shows that a better bound is available in terms of $\mathrm{e}_{\text {conv }}^{\theta}(m)$, Conjecture 133 and Question 136 could be relevant for this purpose.
5.4. Gentle partitions of unity. The following definition describes a numerical parameter that underlies the extension method of [LN05].

Definition 141 (modulus of gentle partition of unity). Suppose that $\left(m, d_{m}\right)$ is a metric space and that $\mathcal{C} \subseteq m$ is nonempty and closed. Define the modulus of gentle partition of unity of $m$ relative to $\mathcal{C}$, denoted $\mathrm{GPU}\left(m, d_{m} ; \mathrm{C}\right)$ or simply $\mathrm{GPU}(m ; \mathrm{C})$ when the metric is clear from the context, to be the infimum over those $g \in(0, \infty]$ such that for every $x \in M$ there is a Borel probability measure $\mu_{x}$ supported on $\mathcal{C}$ with the requirements that if $s \in \mathcal{C}$, then $\mu_{s}=\boldsymbol{\delta}_{s}$, and also for every $x, y \in M$ we have

$$
\int_{\mathbb{C}} d_{m}(s, x) \mathrm{d}\left|\mu_{x}-\mu_{y}\right|(s) \leqslant g d_{m}(x, y) .
$$

The modulus of gentle partitions of unity of $m$, denoted $\operatorname{GPU}\left(m, d_{m}\right)$ or simply $\operatorname{GPU}(m)$ when the metric is clear from the context, is the supremum of $\operatorname{GPU}\left(m, d_{m} ; \mathcal{C}\right)$ over all nonempty closed subsets $\mathrm{C} \subseteq m$.

The nomenclature of Definition 141 is derived from [LN05], though we warn that Definition 141 considers objects that are not identical to those that were introduced in [LN05]. In [LN05] the measures $\left\{\mu_{x}\right\}_{x \in m} M \backslash \mathcal{e}$ were also required to have a Radon-Nikoým derivative with respect to some reference measure $\mu$. This additional requirement arises automatically from the constructions of [LN05] but it is not needed for any of the known applications of gentle partitions of unity, so it is beneficial to remove it altogether. The formal connection between [LN05] and Definition 141 was clarified in [AP20].

In anticipation of the proof of Theorem 65, one can generalize Definition 141 to the case of general profiles, analogously to what we did in Definition 63 .
Definition 142 (gentle partition of unity profile). Suppose that ( $M, d_{m}$ ) is a metric space and that $\mathcal{C} \subseteq m$ is nonempty and closed. A metric $\mathfrak{d}: m \times m \rightarrow[0, \infty)$ is called a gentle partition of unity profile for $\left(m, d_{m}\right)$ relative to $\mathcal{C}$ if for every $x \in M$ there is a Borel probability measure $\mu_{x}$ supported on $\mathcal{C}$ with the requirements that if $s \in \mathcal{C}$, then $\mu_{s}=\boldsymbol{\delta}_{s}$, and also for every $x, y \in M$ we have

$$
\int_{\mathfrak{C}} d_{m}(s, x) \mathrm{d}\left|\mu_{x}-\mu_{y}\right|(s) \leqslant \mathfrak{d}(x, y) .
$$

Ifd is a gentle partition of unity profile for $\left(m, d_{m}\right)$ relative to every closed $\varnothing \neq \mathcal{C} \subseteq m$, then we say that $\mathfrak{d}$ is a gentle partition of unity profile for $(m, d m)$.

Note in passing that if $\mathfrak{d}$ is a gentle partition of unity profile for $\left(m, d_{m}\right)$ relative to $\mathfrak{C}$, then for every $x \in W$ the probability measure $\mu_{x}$ in Definition 142 has finite first moment. Indeed, for any $s_{0} \in \mathcal{C}$,

$$
\begin{equation*}
\int_{\mathrm{e}} d_{m}\left(s_{0}, s\right) \mathrm{d} \mu_{x}(s)=\int_{\mathrm{C}} d_{m}\left(s_{0}, s\right) \mathrm{d}\left(\mu_{x}-\delta_{s_{0}}\right)(s) \leqslant \int_{\mathrm{e}} d_{m}\left(s_{0}, s\right) \mathrm{d}\left|\mu_{x}-\mu_{s_{0}}\right|(s) \leqslant \mathfrak{d}\left(s_{0}, x\right)<\infty, \tag{294}
\end{equation*}
$$

where we used the fact that $\mu_{s_{0}}=\delta_{s_{0}}$, since $s_{0} \in \mathcal{C}$.
Suppose that $\left(m, d_{m}\right)$ is a Polish metric space. The following estimate is implicit in LN05.

$$
\begin{equation*}
\mathrm{e}_{\mathrm{conv}}(m) \leqslant 2 \mathrm{GPU}(m) . \tag{295}
\end{equation*}
$$

In fact, the same reasoning as in LN05 leads to the following more general lemma.
Lemma 143. Suppose that $\left(m, d_{m}\right)$ is a Polish metric space and that $\mathcal{C} \subseteq m$ is nonempty and closed. Assume that $\mathfrak{d}: m \times m \rightarrow[0, \infty)$ is a gentle partition of unity profile for $\left(m, d_{m}\right)$ relative to $\mathcal{C}$. Then, for every Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ and every 1-Lipschitz mapping $f: \mathcal{C} \rightarrow \mathbf{Z}$ there exists

$$
F: m \rightarrow \overline{\operatorname{conv}}(f(\mathrm{C}))
$$

that extends $f$ and satisfies $\|F(x)-F(y)\|_{\mathbf{z}} \leqslant 2 \mathfrak{d}(x, y)$ for every $x, y \in m$.
Proof. Let $\left\{\mu_{x}\right\}_{x \in m}$ be probability measures as in Definition 142. Then, $\left\{\mu_{x}\right\}_{x \in M} \subseteq \mathrm{P}_{1}$ ( C ) by (294). So, by Proposition 135 (with $\varepsilon \equiv 0$ ) it suffices to check that $\mathrm{W}_{1}\left(\mu_{x}, \mu_{y}\right) \leqslant 2 \mathfrak{d}(x, y)$ for every $x, y \in W$. To this end, fix $\eta>0$ and $s_{0} \in \mathcal{C}$ such that $d_{m}\left(x, s_{0}\right) \leqslant d_{m}(x, \mathcal{C})+\eta$. Then,

$$
\forall s \in \mathcal{C}, \quad d_{m}\left(s, s_{0}\right) \leqslant d_{m}(s, x)+d_{m}\left(x, s_{0}\right) \leqslant d_{m}(s, x)+d_{m}(x, \mathcal{C})+\eta \leqslant 2 d_{m}(s, x)+\eta
$$

Consequently, every 1-Lipschitz function $\psi: \mathcal{C} \rightarrow \mathbb{R}$ satisfies

$$
\begin{aligned}
\int_{\mathbb{C}} \psi \mathrm{d} \mu_{x}-\int_{\mathbb{C}} \psi \mathrm{d} \mu_{y}=\int_{\mathbb{C}}\left(\psi(s)-\psi\left(s_{0}\right)\right) \mathrm{d}\left(\mu_{x}-\mu_{y}\right)(s) & \leqslant \int_{\mathbb{C}}\left|\psi(s)-\psi\left(s_{0}\right)\right| \mathrm{d}\left|\mu_{x}-\mu_{y}\right|(s) \\
& \leqslant \int_{\mathbb{C}} d_{m}\left(s, s_{0}\right) \mathrm{d}\left|\mu_{x}-\mu_{y}\right|(s) \leqslant \int_{\mathbb{C}}\left(2 d_{m}(s, x)+\eta\right) \mathrm{d}\left|\mu_{x}-\mu_{y}\right|(s) \leqslant 2 \mathfrak{d}(x, y)+2 \eta .
\end{aligned}
$$

The desired conclusion follows by letting $\eta \rightarrow 0$ and using the Kantorovich-Rubinstein duality (288).
5.5. The multi-scale construction. Suppose that $\left(m, d_{m}\right)$ is a Polish metric space and fix another metric $\mathfrak{d}$ on $m$. In this section we will show that there is a universal constant $\alpha \geqslant 1$ with the following property. Assume that either $\left(m, d_{m}\right)$ is locally compact and $\mathfrak{d}$ is a separation modulus for $\left(m, d_{m}\right)$ per Definition 63, or the assumptions of Theorem 113 are satisfied. We will prove that either of these assumptions implies that $\alpha \mathfrak{d}$ is a gentle partition of unity profile for $\left(m, d_{m}\right)$. By Lemma 143 this gives Theorem 65 and Theorem 113 and will show that in fact these extension results are both convex hull-valued and via a linear extension operator. This also implies that every locally compact metric space $m$ satisfies

$$
\begin{equation*}
\operatorname{GPU}(m) \lesssim \operatorname{SEP}(m) . \tag{296}
\end{equation*}
$$

Remark 144. The bound (296) need not be sharp. Indeed, it was proved in LN05 that if $m$ is finite, then

$$
\begin{equation*}
\operatorname{GPU}(m) \lesssim \frac{\log |m|}{\log \log |m|} \tag{297}
\end{equation*}
$$

However, by Bar96 sometimes SEP $(m) \gtrsim \log |m|$ (and always $\operatorname{SEP}(m) \lesssim \log |m|$ ). A shorter presentation of the proof of (297) can be found in [Nao15], and a different proof of 297) will appear in the forthcoming work [MN21. Also, in the forthcoming work [MNR21] it is proved that 297) is optimal.

The following theorem is a precise formulation of what we will prove in this section.
Theorem 145. Let $(m, d m)$ be a Polish metric space and fix another metric $\mathfrak{d}$ on $m$. Suppose that for every $\Delta>0$ there is a probability space $\left(\Omega_{\Delta}, \mathbf{P r o b}_{\Delta}\right)$ and a sequence of set-valued mappings $\left\{\Gamma_{\Delta}^{k}: \Omega_{\Delta} \rightarrow 2^{m_{j}}\right\}_{k=1}^{\infty}$ such that one of the following two measurability assumptions hold.

- Either $\left(m, d_{m}\right)$ is locally compact and $\Gamma_{\Delta}^{k}$ is strongly measurable for each fixed $k \in \mathbb{N}$ and $\Delta>0$,
- or $\Omega_{\Delta}$ is a Borel subset of some Polish metric space $\mathscr{L}_{\Delta}$ and Prob $_{\Delta}$ is a Borel probability measure supported on $\Omega_{\Delta}$, and $\Gamma_{\Delta}^{k}$ is a standard set-valued mapping for each fixed $k \in \mathbb{N}$ and $\Delta>0$.
Suppose that the following three requirements hold.
(1) $\mathcal{P}_{\Delta}^{\omega}=\left\{\Gamma_{\Delta}^{k}(\omega)\right\}_{k=1}^{\infty}$ is a partition of $m$ for every $\omega \in \Omega_{\Delta}$,
(2) $\operatorname{diam} m\left({ }_{\Delta}^{( }{ }_{\Delta}^{\omega}(x)\right)<\Delta$ for every $x \in M$ and $\omega \in \Omega_{\Delta}$,
(3) $\Delta \operatorname{Prob}_{\Delta}\left[\omega \in \Omega_{\Delta}: \mathcal{P}_{\Delta}^{\omega}(x) \neq \mathcal{P}_{\Delta}^{\omega}(y)\right] \leqslant \mathfrak{d}(x, y)$ for every $x, y \in T$.

Then, $\alpha \mathfrak{d}$ is a gentle partition of unity profile for $\left(m, d_{m}\right)$ for some universal constant $\alpha \in[1, \infty)$.
Suppose from now on that $\mathcal{C}$ is a nonempty closed subset of $m$. We will first set notation and record basic properties of a sequence of bump functions that will be used in the proof of Theorem (145); this part of the discussion is entirely standard and has nothing to do with random partitions.

Fix a 1-Lipschitz function $\psi:[0, \infty) \rightarrow[0, \infty)$ such that $\operatorname{supp}(\psi) \subseteq[1,4]$ and $\psi(t)=1$ for every $t \in[2,3]$ (these requirements uniquely determine $\psi$, which is piecewise linear). Define for each $n \in \mathbb{Z}$,

$$
\forall x \in m, \quad \phi_{n}(x)=\phi_{n}^{\text {е. }}(x) \stackrel{\text { def }}{=} \psi\left(2^{-n} d_{m}(x, \mathcal{C})\right) .
$$

Then $\left\|\phi_{n}\right\|_{\text {Lip }(m)} \leqslant 2^{-n}$ and if $\phi_{n}(x) \neq 0$ then necessarily $2^{n} \leqslant d_{m}(x, \mathcal{C}) \leqslant 2^{n+2}$. We also denote

$$
\forall x \in m, \quad \Phi(x)=\Phi^{\mathcal{C}}(x) \stackrel{\operatorname{def}}{=} \sum_{m \in \mathbb{Z}} \phi_{n}(x) .
$$

For each $x \in M$, at most two summands in the sum that defines $\Phi(x)$ do not vanish. If $x \in m \backslash \mathcal{C}$, then since $\mathcal{C}$ is closed we have $d_{m}(x, \mathcal{C})>0$, and therefore there is $n \in \mathbb{Z}$ for which $2^{n} \leqslant d_{m}(x, \mathcal{C})<2^{n+1}$. For this value of $n$ we have $\phi_{n}(x)=1$, so $\Phi(x) \geqslant 1$ for every $x \in m \backslash \mathcal{C}$. Finally, for each $n \in \mathbb{Z}$ define

$$
\forall x \in M, \quad \lambda_{n}(x)=\lambda_{n}^{\mathcal{C}}(x) \stackrel{\operatorname{def}}{=} \begin{cases}\frac{\phi_{n}(x)}{\Phi(x)} & \text { if } x \in M \backslash \mathcal{C}, \\ 0 & \text { if } x \in \mathcal{C} .\end{cases}
$$

By design, $\sum_{n \in \mathbb{Z}} \lambda_{n}(x)=1$ for every $x \in M \backslash \mathcal{C}$. Further properties of these bump functions are recorded in the following basic lemma, for ease of later reference.

Lemma 146. Suppose that $x, y \in M$ satisfy $d_{m}(x, \mathcal{C}) \geqslant d_{m}(y, \mathcal{C})>d_{m}(x, y)$. Then for every $n \in \mathbb{Z}$,

$$
\begin{equation*}
\frac{2^{n}}{d_{m}(y, \mathrm{C})} \notin\left(\frac{1}{4}, 2\right) \Longrightarrow \phi_{n}(x)=\phi_{n}(y)=\lambda_{n}(x)=\lambda_{n}(y)=0, \tag{298}
\end{equation*}
$$

and

$$
\begin{equation*}
2^{n-1}<d_{m}(y, \mathcal{C})<2^{n+2} \Longrightarrow\left|\lambda_{n}(x)-\lambda_{n}(y)\right| \lesssim \frac{d_{m}(x, y)}{d_{m}(y, \mathcal{C})} . \tag{299}
\end{equation*}
$$

Proof. Our assumption implies that $d_{m}(x, \mathcal{C}), d_{m}(y, \mathcal{C})>0$, so $x, y \in m \backslash \mathcal{C}$. To prove 298, suppose first that $2^{n} \geqslant 2 d_{m}(y, \mathcal{C})$. Then, $\phi_{n}(y)=\lambda_{n}(y)=0 \operatorname{since} \operatorname{supp}(\psi) \subseteq[1,4]$ and $2^{-n} d_{m}(y, \mathrm{C}) \leqslant 1$. Also, $d_{m}(x, \mathcal{C}) \leqslant d_{m}(x, y)+d_{m}(y, \mathcal{C})<2 d_{m}(y, \mathcal{C}) \leqslant 2^{n}$, so $2^{-n} d_{m}(x, \mathcal{C}) \leqslant 1$ and hence $\phi_{n}(x)=\lambda_{n}(x)=0$. The remaining case of (298) is when $d_{m}(y, \mathcal{C}) \geqslant 2^{n+2}$. Then, $2^{-n} d_{m}(x, \mathcal{C}) \geqslant 2^{-n} d_{m}(y, \mathcal{C}) \geqslant 4$ and therefore $\left\{2^{-n} d_{m}(x, \mathcal{C}), 2^{-n} d_{m}(y, \mathcal{C})\right\} \cap \operatorname{supp}(\psi)=\varnothing$. Consequently, $\phi_{n}(x)=\phi_{n}(y)=\lambda_{n}(x)=\lambda_{n}(y)=0$.

To prove (299), assume that $2^{n-1}<d_{m}(y, \mathcal{C})<2^{n+2}$. Recalling that (point-wise) on $m \backslash \mathcal{C}$ we have $\lambda_{n}=\phi_{n} / \Phi$ for all $n \in \mathbb{Z}$ and $\Phi \geqslant 1$, and moreover $\left\|\phi_{n}\right\|_{\text {Lip }(m)} \leqslant 2^{-n}$, we conclude as follows.

$$
\begin{aligned}
\left|\lambda_{n}(x)-\lambda_{n}(y)\right| & \leqslant \frac{\left|\phi_{n}(x)-\phi_{n}(y)\right|}{\Phi(x)}+\frac{\phi_{n}(y)}{\Phi(x) \Phi(y)}|\Phi(y)-\Phi(x)| \\
& \leqslant 2^{-n} d_{m}(x, y)+\sum_{n \in \mathbb{Z}}\left|\phi_{n}(x)-\phi_{n}(y)\right| \\
& \stackrel{298}{\leqslant} 2^{-n} d_{m}(x, y)+\sum_{\substack{n \in \mathbb{Z} \\
2^{n-1}<d_{m}(y, \mathrm{C})<2^{n+2}}} 2^{-n} d_{m}(x, y)=\frac{d_{m}(x, y)}{d_{m}(y, \mathcal{C})} .
\end{aligned}
$$

The interaction between $\left\{\lambda_{n}\right\}_{n \in \mathbb{Z}}$ and the random partitions of Theorem 145 is the content of the following lemma. Note that by reasoning as in (95), the metric $\mathfrak{d}$ in Theorem 145 must satisfy

$$
\forall x, y \in M, \quad \mathfrak{d}(x, y) \geqslant d_{m}(x, y) .
$$

Lemma 147. In the setting of Theorem 145, if $x \in T \backslash \mathcal{C}$ and $y \in M \backslash\{x\} \operatorname{satisfy} d_{m}(x, \mathcal{C}) \geqslant d_{m}(y, \mathcal{C})$, then

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}} \sum_{k=1}^{\infty} \int_{\Omega_{2^{n}}}\left|\lambda_{n}(x) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(x)-\lambda_{n}(y) \mathbf{1}_{\Gamma_{2^{n}}(\omega)}(y)\right| \operatorname{dProb}_{2^{n}}(\omega) \lesssim \frac{\mathfrak{d}(x, y)}{d_{m}(y, \mathcal{C})+d_{m}(x, y)} . \tag{300}
\end{equation*}
$$

Proof. As $\sum_{n \in \mathbb{Z}} \lambda_{n}(x)=\sum_{n \in \mathbb{Z}} \lambda_{n}(y)=1$ and also $\sum_{k=1}^{\infty} \mathbf{1}_{\Gamma_{2^{n}}(\omega)}(x)=\sum_{k=1}^{\infty} \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(y)=1$ for each $n \in \mathbb{Z}$ and $\omega \in \Omega_{2^{n}}$, the left hand side of (300) is at most 2. Since $\mathfrak{d}(x, y) \geqslant d_{m}(x, y)$, it follows that (300) holds if $d_{m}(y, \mathrm{C}) \leqslant d_{m}(x, y)$. So, we will assume in the rest of the proof of Lemma 147 that $d_{m}(x, y)<d_{m}(y, \mathrm{C})$ (thus, in particular, $y \in m \backslash \mathcal{C}$ ), in which case the right-hand side of (300) becomes at least a universal constant multiple of the quantity $\mathfrak{d}(x, y) / d_{m}(y, \mathcal{C})$.

We claim that for every $n \in \mathbb{Z}$ the following inequality holds for every $\omega \in \Omega_{2^{n}}$.

$$
\begin{equation*}
\sum_{k=1}^{\infty}\left|\lambda_{n}(x) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(x)-\lambda_{n}(y) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(y)\right| \lesssim\left(2^{-n} d m(x, y)+\mathbf{1}_{\left\{P_{2^{m}}^{\omega}(x) \neq \mathcal{P}_{2^{n}}^{\omega}(y)\right\}}\right) \mathbf{1}_{\left\{\frac{1}{4}<\frac{2^{n}}{d_{m}(y, C)}<2\right\}} . \tag{301}
\end{equation*}
$$

Assuming (301), we conclude the proof of (300) in the remaining case $d_{m}(x, y)<d_{m}(y, \mathcal{C})$ as follows.

$$
\begin{aligned}
\sum_{n \in \mathbb{Z}} \sum_{k=1}^{\infty} & \int_{\Omega_{2^{n}}}\left|\lambda_{n}(x) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(x)-\lambda_{n}(y) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(y)\right| \mathrm{d}_{\operatorname{Pro}}^{2^{n}} \\
& \vdots \sum_{\substack{n \in \mathbb{Z} \\
2^{n-1}<d_{m}(y, \mathcal{C})<2^{n+2}}}\left(2^{-n} d_{m}(x, y)+\operatorname{Prob}_{2^{n}}\left[\left\{\omega \in \Omega_{2^{n}}: \mathcal{P}_{2^{n}}^{\omega}(x) \neq \mathcal{P}_{2^{n}}^{\omega}(y)\right\}\right]\right) \\
& \lesssim \sum_{\substack{n \in \mathbb{Z} \\
2^{m n-1}<d_{m}(y, \mathrm{C})<2^{n+2}}} 2^{-n}\left(d_{m}(x, y)+\mathfrak{d}(x, y)\right)=\frac{\mathfrak{d}(x, y)}{d_{m}(y, \mathcal{C})}=\frac{\mathfrak{d}(x, y)}{d_{m}(y, \mathcal{C})+d_{m}(x, y)},
\end{aligned}
$$

where the first step uses (301), the second step is where we used condition (3) of Theorem 145 , the penultimate step uses $\mathfrak{d}(x, y) \geqslant d_{m}(x, y)$, and in the final step uses the assumption $d_{m}(x, y)<d_{m}(y, \mathcal{C})$.

It remains to establish (301. By Lemma 146 . if it is not the case that $2^{n-1}<d_{m}(y, \mathcal{C})<2^{n+2}$, then $\lambda_{n}(x)=\lambda_{n}(y)=0$, so both sides of (301) vanish. We may therefore assume that $2^{n-1}<d m(y, \mathcal{C})<2^{n+2}$. Under this assumption, if $\mathcal{P}_{2^{n}}^{\omega}(x) \neq \mathcal{P}_{2^{n}}^{\omega}(y)$, then the right-hand side of (301) is at least 1 , while the lefthand side of (301) consists of a sum of two numbers, each of which is at most 1 . It therefore remains to establish (301) when $\mathcal{P}_{2^{n}}^{\omega}(x)=\mathcal{P}_{2^{n}}^{\omega}(y)$ (and still $2^{n-1}<d m(y, \mathcal{C})<2^{n+2}$ ). In this case, 301) becomes the inequality $\left|\lambda_{2^{n}}(x)-\lambda_{2^{n}}(y)\right| \leqslant d_{m}(x, y) / d_{m}(y, \mathcal{C})$, which we proved in Lemma 146 .

Proof of Theorem 145, By Lemma 114 and Corollary 117, for every $\Delta>0$ there exists a Prob ${ }_{\Delta}$-to-Borel measurable mapping $\gamma_{\Delta}^{k}: \Omega_{m} \rightarrow \mathcal{C}$ such that

$$
\begin{equation*}
\forall \omega \in \Omega_{\Delta}, \quad \Gamma_{\Delta}^{k}(\omega) \neq \varnothing \Longrightarrow d m\left(\gamma_{\Delta}^{k}(\omega), \Gamma_{\Delta}^{k}(\omega)\right) \leqslant d_{m}\left(\mathcal{C}, \Gamma_{\Delta}^{k}(\omega)\right)+\Delta \tag{302}
\end{equation*}
$$

(In fact, in the locally compact setting of Theorem 145 , the use of Lemma 114 shows that the additive $\Delta$ term in the right hand side of (302) can be removed).

For every $x \in M \backslash \mathcal{C}$ define a Borel measure $\mu_{x}$ supported on $\mathcal{C}$ by

$$
\begin{equation*}
\mu_{x} \stackrel{\text { def }}{=} \sum_{n \in \mathbb{Z}} \sum_{k=1}^{\infty} \lambda_{n}(x)\left(\gamma_{2^{n}}^{k}\right)_{\#}\left(\operatorname{Prob}_{2^{n}} L_{\left\{\omega \in \Omega_{2^{n}}: x \in \Gamma_{2^{n}}^{k}(\omega)\right\}}\right) \tag{303}
\end{equation*}
$$

In other words, for every Borel-measurable mapping $h: \mathcal{C} \rightarrow[0, \infty)$ we have

$$
\begin{equation*}
\int_{\mathcal{C}} h(s) \mathrm{d} \mu_{x}(s)=\sum_{n \in \mathbb{Z}} \sum_{k=1}^{\infty} \lambda_{n}(x) \int_{\left\{\omega \in \Omega_{2^{n}}: x \in \Gamma_{2^{n}}^{k}(\omega)\right\}} h\left(\gamma_{2^{n}}^{k}(\omega)\right) \mathrm{d}_{\mathbf{P r o b}}^{2^{n}}{ }^{(\omega)} \tag{304}
\end{equation*}
$$

Since $\mathcal{P}_{2^{n}}^{\omega}$ is a partition of $X$ for every $n \in \mathbb{Z}$ and $\omega \in \Omega_{2^{n}}$, the special case $h=\mathbf{1}_{\mathcal{C}}$ of (304) implies that

$$
\begin{aligned}
\mu_{x}(\mathcal{C}) & =\sum_{n \in \mathbb{Z}} \sum_{k=1}^{\infty} \lambda_{n}(x) \operatorname{Prob}_{2^{n}}\left[\left\{\omega \in \Omega_{2^{n}}: x \in \Gamma_{2^{n}}^{k}(\omega)\right\}\right] \\
& =\sum_{n \in \mathbb{Z}} \lambda_{n}(x) \operatorname{Prob}_{2^{n}}\left[\left\{\omega \in \Omega_{2^{n}}: x \in \bigcup_{k=1}^{\infty} \Gamma_{2^{n}}^{k}(\omega)\right\}\right]=\sum_{n \in \mathbb{Z}} \lambda_{n}(x)=1
\end{aligned}
$$

Thus $\mu_{x}$ is a probability measure. Consequently, if we also denote $\mu_{s}=\boldsymbol{\delta}_{s}$ for every $s \in \mathcal{C}$, then the proof of Theorem 145 will be complete if we show that

$$
\begin{equation*}
\forall x, y \in M, \quad \int_{\mathcal{C}} d m(s, x) \mathrm{d}\left|\mu_{x}-\mu_{y}\right|(s) \lesssim \mathfrak{d}(x, y) \tag{305}
\end{equation*}
$$

It suffices to prove (305) when $x, y \in M$ are distinct and $\{x, y\} \nsubseteq \mathcal{C}$. Indeed, if $\{x, y\} \subseteq \mathcal{C}$ then $\mu_{x}=\boldsymbol{\delta}_{x}$ and $\mu_{y}=\boldsymbol{\delta}_{y}$, so the left hand side of (305) is equal to $d_{m}(x, y)$, which is at most $\mathfrak{d}(x, y)$. Hence, in the rest of the proof of Theorem 145 we will assume without loss of generality that $x \in M \backslash \mathcal{C}$ and $d_{m}(x, \mathcal{C}) \geqslant d_{m}(y, \mathcal{C})$.

We claim that the left hand side of 305) can be bounded from above as follows.

$$
\begin{align*}
& \int_{\mathcal{C}} d m(s, x) \mathrm{d}\left|\mu_{x}-\mu_{y}\right|(s) \\
& \quad \leqslant d_{m}(x, y)+\sum_{n \in \mathbb{Z}} \sum_{k=1}^{\infty} \int_{\Omega_{2^{n}}} d_{m}\left(\gamma_{2^{n}}^{k}(\omega), x\right)\left|\lambda_{n}(x) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(x)-\lambda_{2^{n}}(y) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(y)\right| \mathrm{d} \operatorname{Prob}_{2^{n}}(\omega) \tag{306}
\end{align*}
$$

Indeed, if $x, y \in M \backslash \mathcal{C}$, then $\mu_{x}, \mu_{y}$ are defined according to 303, so that

$$
\begin{aligned}
& \int_{\mathcal{C}} d m(s, x) \mathrm{d}\left|\mu_{x}-\mu_{y}\right|(s) \\
& \quad \leqslant \sum_{n \in \mathbb{Z}} \sum_{k=1}^{\infty} \int_{\mathcal{C}} d_{m}(s, x) \mathrm{d}\left(\left(\gamma_{2^{n}}^{k}\right)_{\#} \mid \lambda_{n}(x) \operatorname{Prob}_{2^{n}}\left\lfloor_{\left\{\omega \in \Omega_{2^{n}}: x \in \Gamma_{2^{n}}^{k}(\omega)\right\}}-\lambda_{n}(y) \operatorname{Prob}_{2^{n}}\left\lfloor_{\left\{\omega \in \Omega_{2^{n}}: y \in \Gamma_{2^{n}}^{k}(\omega)\right\}} \mid\right)(s)\right.\right. \\
& \quad=\sum_{n \in \mathbb{Z}} \sum_{k=1}^{\infty} \int_{\Omega_{2^{n}}} d_{m}\left(\gamma_{2^{n}}^{k}(\omega), x\right)\left|\lambda_{n}(x) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(x)-\lambda_{n}(y) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(y)\right| \operatorname{dProb}_{2^{n}}(\omega)
\end{aligned}
$$

thus establishing (306) in this case. The remaining case is when $x \in M \backslash \mathcal{C}$ and $y \in \mathcal{C}$, so that $\mu_{x}$ is given in (303) and $\mu_{y}=\boldsymbol{\delta}_{y}$. We can then use the following (crude) estimate.

$$
\begin{align*}
\int_{\mathcal{C}} d_{m}(s, x) \mathrm{d}\left|\mu_{x}-\mu_{y}\right|(s) & \leqslant \int_{\mathcal{C}} d m(s, x) \mathrm{d} \mu_{y}(s)+\int_{\mathcal{C}} d m(s, x) \mathrm{d} \mu_{x}(s) \\
& =d_{m}(x, y)+\sum_{n \in \mathbb{Z}} \sum_{k=1}^{\infty} \int_{\Omega_{2^{n}}} d_{m}\left(\gamma_{2^{n}}^{k}(\omega), x\right) \lambda_{n}(x) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(x) \mathrm{d} \mathbf{P r o b}_{2^{n}}(\omega) \tag{307}
\end{align*}
$$

It remains to observe that because $y \in \mathcal{C}$ we have $\lambda_{n}(y)=0$ for all $n \in \mathbb{Z}$ and therefore the right hand side of (307) coincides with the right hand side of (306).

Next, we claim that for every $(n, k) \in \mathbb{Z} \times \mathbb{N}$ and every $\omega \in \Omega_{2^{n}}$ we have

$$
\begin{align*}
d_{m}\left(\gamma_{2^{n}}^{k}(\omega), x\right) \mid \lambda_{n}(x) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(x) & -\lambda_{n}(y) \mathbf{1}_{\Gamma_{2^{n}(\omega)}^{k}}(y) \mid  \tag{308}\\
& \lesssim\left(d_{m}(y, \mathrm{C})+d_{m}(x, y)\right)\left|\lambda_{n}(x) \mathbf{1}_{\Gamma_{2^{n}}^{k}(\omega)}(x)-\lambda_{n}(y) \mathbf{1}_{\Gamma_{2^{n}( }^{k}(\omega)}(y)\right|
\end{align*}
$$

By a substitution of the point-wise estimate (308) into 306) and using $d_{m}(x, y) \leqslant \mathfrak{d}(x, y)$ the desired estimate (305) follows from Lemma 147, thus completing the proof of Theorem 145 ,

To verify (308), note first that both sides of (308) vanish unless $x \in \Gamma_{2^{n}}^{k}(\omega)$ or $y \in \Gamma_{2^{n}}^{k}(\omega)$ and also, due to Lemma $146,2^{n-1}<d_{m}(y, \mathcal{C})<2^{n+2}$. So, assume from now on that

$$
\begin{equation*}
\{x, y\} \cap \Gamma_{2^{n}}^{k}(\omega) \neq \varnothing \quad \text { and } \quad 2^{n-1}<d_{m}(y, \mathcal{C})<2^{n+2} \tag{309}
\end{equation*}
$$

Our goal (308) then becomes to deduce that

$$
\begin{equation*}
d_{m}\left(\gamma_{2^{n}}^{k}(\omega), x\right) \lesssim d_{m}(y, \mathcal{C})+d_{m}(x, y) \tag{310}
\end{equation*}
$$

Choose a point $z \in \Gamma_{m}^{k}(\omega)$ such that

$$
\begin{equation*}
d_{m}\left(\gamma_{2^{n}}^{k}(\omega), z\right) \leqslant d_{m}\left(\gamma_{2^{n}}^{k}(\omega), \Gamma_{2^{n}}^{k}(\omega)\right)+2^{n} \stackrel{302}{=} d_{m}\left(\mathcal{C}, \Gamma_{2^{n}}^{k}(\omega)\right)+2^{n+1} \stackrel{309}{=} d_{m}\left(\mathcal{C}, \Gamma_{2^{n}}^{k}(\omega)\right)+d_{m}(y, \mathcal{C}) \tag{311}
\end{equation*}
$$

If $x \in \Gamma_{2^{n}}^{k}(\omega)$, then
$d_{m}\left(\mathcal{C}, \Gamma_{2^{n}}^{k}(\omega)\right) \leqslant d_{m}(x, \mathcal{C}) \leqslant d_{m}(x, y)+d_{m}(y, \mathcal{C}) \quad$ and $\quad d_{m}(x, z) \leqslant \operatorname{diam}_{m}\left(\Gamma_{2^{n}}^{k}(\omega)\right) \leqslant 2^{n} \stackrel{309}{=} d_{m}(y, \mathcal{C})$.
By combining these two estimates with (311) and the triangle inequality, we see that

$$
d_{m}\left(\gamma_{2^{n}}^{k}(\omega), x\right) \leqslant d_{m}\left(\gamma_{2^{n}}^{k}(\omega), z\right)+d_{m}(z, x) \lesssim d_{m}(x, y)+d_{m}(y, \mathcal{C})
$$

Hence, the desired estimate (310) holds when $x \in \Gamma_{2^{n}}^{k}(\omega)$.
It remains to check (310) when $y \in \Gamma_{2^{n}}^{k}(\omega)$, in which case we proceed similarly by noting that now

$$
d_{m}\left(\mathcal{C}, \Gamma_{2^{n}}^{k}(\omega)\right) \leqslant d_{m}(y, \mathcal{C}) \quad \text { and } \quad d_{m}(y, z) \leqslant \operatorname{diam}_{m}\left(\Gamma_{2^{n}}^{k}(\omega)\right) \leqslant 2^{n \stackrel{309}{=}} d_{m}(y, \mathcal{C})
$$

By combining these two estimates with (311) and the triangle inequality, we conclude that

$$
d_{m}\left(\gamma_{2^{n}}^{k}(\omega), x\right) \leqslant d_{m}\left(\gamma_{2^{n}}^{k}(\omega), z\right)+d_{m}(z, y)+d_{m}(y, x) \lesssim d_{m}(y, \mathcal{C})+d_{m}(x, y)
$$

## 6. Volume computations

In this section we will prove volume estimates that occur in our bounds on the separation modulus.
6.1. Direct sums. Fix $n \in \mathbb{N}$ and a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$. Throughout what follows, the (normalized) cone measure [GM87] on $\partial B_{\mathbf{X}}$ will be denoted $\kappa_{\mathbf{X}}$. Thus, for every measurable $A \subseteq \partial B_{\mathbf{X}}$,

$$
\begin{equation*}
\kappa_{\mathbf{X}}(A) \stackrel{\text { def }}{=} \frac{\operatorname{vol}_{n}([0,1] A)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}=\frac{\operatorname{vol}_{n}(\{s v:(s, v) \in[0,1] \times A\})}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} . \tag{312}
\end{equation*}
$$

The probability measure $\kappa_{\mathbf{X}}$ is characterized by the following "generalized polar coordinates" identity, which holds for every $f \in L_{1}\left(\mathbb{R}^{n}\right)$; see e.g. [NR03, Proposition 1].

$$
\begin{equation*}
\int_{\mathbb{R}^{n}} f(x) \mathrm{d} x=n \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) \int_{0}^{\infty} r^{n-1}\left(\int_{\partial B_{\mathbf{X}}} f(r \theta) \mathrm{d} \kappa_{\mathbf{X}}(\theta)\right) \mathrm{d} r \tag{313}
\end{equation*}
$$

As a quick application of (313), we will next record for ease of later reference the following computation of the volume of the unit ball of an $\ell_{p}$ direct sum of normed spaces.
Lemma 148. Fix $n, m_{1}, \ldots, m_{n} \in \mathbb{N}$ and normed spaces $\left\{\mathbf{X}_{j}=\left(\mathbb{R}^{m_{1}},\|\cdot\|_{\mathbf{X}_{m_{j}}}\right)\right\}_{j=1}^{n}$. Then

$$
\begin{equation*}
\forall p \in[1, \infty], \quad \operatorname{vol}_{m_{1}+\ldots+m_{n}}\left(B \mathbf{X}_{1} \oplus_{p} \ldots \oplus_{p} \mathbf{X}_{n}\right)=\frac{\prod_{j=1}^{n} \Gamma\left(1+\frac{m_{j}}{p}\right) \operatorname{vol}_{m_{j}}\left(B_{\mathbf{X}_{j}}\right)}{\Gamma\left(1+\frac{m_{1}+\ldots+m_{n}}{p}\right)} \tag{314}
\end{equation*}
$$

Proof. This follows by induction on $n$ from the following identity (direct application of Fubini), which holds for every $a, b \in \mathbb{N}$ and any two normed spaces $\mathbf{X}=\left(\mathbb{R}^{a},\|\cdot\|_{\mathbf{X}}\right)$ and $\mathbf{Y}=\left(\mathbb{R}^{b},\|\cdot\|_{\mathbf{Y}}\right)$.

$$
\begin{aligned}
\operatorname{vol}_{a+b}\left(B_{\mathbf{X} \oplus_{p} \mathbf{Y}}\right)= & \int_{B_{\mathbf{X}}} \operatorname{vol}_{b}\left(\left(1-\|x\|_{\mathbf{X}}^{p}\right)^{\frac{1}{p}} B_{\mathbf{Y}}\right) \mathrm{d} x=\operatorname{vol}_{b}\left(B_{\mathbf{Y}}\right) \int_{B_{\mathbf{X}}}\left(1-\|x\|_{\mathbf{X}}^{p}\right)^{\frac{b}{p}} \mathrm{~d} x \\
& \stackrel{(313)}{=} \operatorname{vol}_{a}\left(B_{\mathbf{X}}\right) \operatorname{vol}_{b}\left(B_{\mathbf{Y}}\right) \int_{0}^{1} a r^{a-1}\left(1-r^{p}\right)^{\frac{b}{p}} \mathrm{~d} r=\operatorname{vol}_{a}\left(B_{\mathbf{X}}\right) \operatorname{vol}_{b}\left(B_{\mathbf{Y}}\right) \frac{\Gamma\left(1+\frac{b}{p}\right) \Gamma\left(1+\frac{a}{p}\right)}{\Gamma\left(1+\frac{a+b}{p}\right)} .
\end{aligned}
$$

By Lemma 148 for every $m \in \mathbb{N}$, every normed space $\mathbf{X}=\left(\mathbb{R}^{m},\|\cdot\|_{\mathbf{X}}\right)$ satisfies

$$
\begin{equation*}
\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}(\mathbf{X})}\right)=\frac{\Gamma\left(1+\frac{m}{p}\right)^{n}}{\Gamma\left(1+\frac{n m}{p}\right)} \operatorname{vol}_{m}\left(B_{\mathbf{X}}\right)^{n} \quad \text { and hence } \quad \operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}(\mathbf{X})}\right)^{\frac{1}{n m}} \simeq \frac{\operatorname{vol}_{m}\left(B_{\mathbf{X}}\right)^{\frac{1}{m}}}{n^{\frac{1}{p}}} \tag{315}
\end{equation*}
$$

In particular, for every $m, n \in \mathbb{N}$ and $1 \leqslant p, q \leqslant \infty$ we have

$$
\begin{equation*}
\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}\left(\ell_{q}^{m}\right)}\right)=\frac{2^{n m} \Gamma\left(1+\frac{1}{q}\right)^{n m} \Gamma\left(1+\frac{m}{p}\right)^{n}}{\Gamma\left(1+\frac{m}{q}\right)^{n} \Gamma\left(1+\frac{n m}{p}\right)} \quad \text { and hence } \quad \operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}\left(\ell_{q}^{m}\right)}\right)^{\frac{1}{n m}}=\frac{1}{n^{\frac{1}{p}} m^{\frac{1}{q}}} \tag{316}
\end{equation*}
$$

The following simple lemma records an extension of the second part of 315) to $m$-fold iterations of the operation $\mathbf{X} \mapsto \ell_{p}^{n}(\mathbf{X})$, i.e., to spaces of the form

$$
\ell_{p_{m}}^{n_{m}}\left(\ell_{p_{m-1}}^{n_{m-1}}\left(\cdots \ell_{p_{1}}^{n_{1}}(\mathbf{X}) \cdots\right)\right)
$$

the main point for us here is that the implicit constants remain bounded as $m \rightarrow \infty$.
Lemma 149. Fix $\left\{n_{k}\right\}_{k=0}^{\infty} \subseteq \mathbb{N}$ and $\left\{p_{k}\right\}_{k=1}^{\infty} \subseteq[1, \infty]$. Let $\mathbf{X}=\left(\mathbb{R}^{n_{0}},\|\cdot\|_{\mathbf{X}}\right)$ be a normed space and define
$\forall k \in \mathbb{N} \cup\{0\}, \quad \mathbf{X}_{k+1}=\ell_{p_{k}}^{n_{k}}\left(\mathbf{X}_{k}\right), \quad$ where $\quad \mathbf{X}_{0}=\mathbf{X}$.
Then, for every $m \in \mathbb{N}$ we have

$$
\operatorname{vol}_{n_{0} \cdots n_{m}}\left(B \mathbf{X}_{m}\right)^{\frac{1}{n_{0} \cdots n_{k}}}=\frac{\operatorname{vol}_{n_{0}}\left(B_{\mathbf{X}}\right)^{\frac{1}{n_{0}}}}{\prod_{k=1}^{m} n_{k}^{\frac{1}{p_{k}}}}
$$

Proof. With the convention that an empty product equals 1, by applying (315) inductively we see that

$$
\operatorname{vol}_{n_{0} \cdots n_{m}}\left(B_{\mathbf{X}_{m}}\right)=\operatorname{vol}_{n_{0}}\left(B_{\mathbf{X}}\right)^{n_{1} \cdots n_{m}} \prod_{k=1}^{m} \frac{\Gamma\left(1+\frac{n_{0} \cdots n_{k-1}}{p_{k}}\right)^{n_{k} \cdots n_{m}}}{\Gamma\left(1+\frac{n_{0} \cdots n_{k}}{p_{k}}\right)^{n_{k+1} \cdots n_{m}}} .
$$

Hence,

$$
\begin{equation*}
\frac{\operatorname{vol}_{n_{0} \cdots n_{m}}\left(B_{\mathbf{X}_{m}}\right)^{\frac{1}{n_{0} \cdots n_{k}}} \prod_{k=1}^{m} n_{k}^{\frac{1}{p_{k}}}}{\operatorname{vol}_{n_{0}}\left(B_{\mathbf{X}}\right)^{\frac{1}{n_{0}}}}=\prod_{k=1}^{m} \frac{\Gamma\left(1+\frac{n_{0} \cdots n_{k-1}}{p_{k}}\right)^{\frac{1}{n_{0} \cdots n_{k-1}}}}{\Gamma\left(1+\frac{n_{0} \cdots n_{k}}{p_{k}}\right)^{\frac{1}{n_{0} \cdots n_{k}}}} n_{k}^{\frac{1}{p_{k}}}=\prod_{k=1}^{m} f_{n_{0} \cdots n_{k-1}, n_{k}}\left(\frac{1}{p_{k}}\right), \tag{317}
\end{equation*}
$$

where for $u, v, t>0$ we denote

$$
f_{u, v}(t) \stackrel{\text { def }}{=} \frac{\Gamma(1+u t)^{\frac{1}{u}}}{\Gamma(1+u v t)^{\frac{1}{u v}}} v^{t} .
$$

Since $(\log \Gamma(z))^{\prime}=\int_{0}^{\infty} \frac{s e^{-z s}}{1-e^{-s}} \mathrm{~d} s$ for $z>0$ (see e.g. [WW62, Chapter XII]), if $u, t>0$ and $v \geqslant 1$, then

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \log f_{u, v}(t)=\log v+\int_{0}^{\infty}\left(e^{-u t s}-e^{-u v t s}\right) \frac{s e^{-s}}{1-e^{-s}} \mathrm{~d} s \geqslant 0 .
$$

Thus, $f_{u, v}$ is increasing on $[0, \infty)$, and therefore we get from (317) that

$$
1=\prod_{k=1}^{m} f_{n_{0} \cdots n_{k-1}, n_{k}}(0) \leqslant \frac{\operatorname{vol}_{n_{0} \cdots n_{m}}\left(B_{\mathbf{X}_{m}}\right)^{\frac{1}{n_{0} \cdots n_{k}}} \prod_{k=1}^{m} n_{k}^{\frac{1}{p_{k}}}}{\operatorname{vol}_{n_{0}}\left(B_{\mathbf{X}}\right)^{\frac{1}{n_{0}}}} \leqslant \prod_{k=1}^{m} f_{n_{0} \cdots n_{k-1}, n_{k}}(1)=\frac{\left(n_{0}!\right)^{\frac{1}{n_{0}}} n_{1} \cdots n_{m}}{\left(\left(n_{0} \cdots n_{m}\right)!\right)^{\frac{1}{n_{0} \cdots n_{m}}}} \leqslant e .
$$

The first part of Lemma 150 below is a restatement of Lemma 36 from the Introduction. Qualitatively, it shows that the class of spaces for which Conjecture 10 holds is closed under unconditional composition, namely, norms of the form (318) below. The second part of Lemma 150 is further information that pertains to Conjecture 48 , i.e., to the symmetric version of the weak reverse isoperimetric conjecture, for which we want the operator $S$ to be the identity mapping (i.e., weak reverse isoperimetry holds without the need to first change the "position" of the given normed space).
Lemma 150. Fix $n, m_{1}, \ldots, m_{n} \in \mathbb{N}$. Let $\mathbf{X}_{1}=\left(\mathbb{R}^{m_{1}},\|\cdot\| \mathbf{x}_{1}\right), \ldots, \mathbf{X}_{n}=\left(\mathbb{R}^{m_{n}},\|\cdot\|_{\mathbf{x}_{n}}\right)$ be normed spaces. Also, let $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ be an unconditional normed space. Define a normed space $\mathbf{X}=\left(\mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}},\|\cdot\|_{\mathbf{X}}\right)$ by

$$
\begin{equation*}
\forall x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}}, \quad\|x\|_{\mathbf{X}} \stackrel{\text { def }}{=}\left\|\left(\left\|x_{1}\right\|_{\mathbf{x}_{1}}, \ldots,\left\|x_{n}\right\|_{\mathbf{x}_{n}}\right)\right\|_{\mathbf{E}} \tag{318}
\end{equation*}
$$

Then, Conjecture 10 (equivalently, Conjecture 35) holds for $\mathbf{X}$ if it holds for $\mathbf{X}_{1}, \ldots, \mathbf{X}_{n}$.
More precisely, suppose that there exist $\alpha>0$, linear transformations $S_{1} \in \mathrm{SL}_{m_{1}}(\mathbb{R}), \ldots, S_{n} \in \mathrm{SL}_{m_{n}}(\mathbb{R})$, and normed spaces $\mathbf{Y}_{1}=\left(\mathbb{R}^{m_{1}},\|\cdot\|_{\mathbf{Y}_{1}}\right), \ldots, \mathbf{Y}_{n}=\left(\mathbb{R}^{m_{n}},\|\cdot\|_{\mathbf{Y}_{n}}\right)$ such that

$$
\begin{equation*}
\forall k \in\{1, \ldots, n\}, \quad B_{\mathbf{Y}_{k}} \subseteq S_{k} B_{\mathbf{X}_{k}} \quad \text { and } \quad \frac{\operatorname{iq}\left(B_{\mathbf{Y}_{k}}\right)}{\sqrt{m_{k}}}\left(\frac{\operatorname{vol}_{m_{k}}\left(B_{\mathbf{X}_{k}}\right)}{\operatorname{vol}_{m_{k}}\left(B_{\mathbf{Y}_{k}}\right)}\right)^{\frac{1}{m_{k}}} \leqslant \alpha . \tag{319}
\end{equation*}
$$

Then, there exist a normed space $\mathbf{Y}=\left(\mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}},\|\cdot\|_{\mathbf{X}}\right)$ and $S \in S L\left(\mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}}\right)$ such that

$$
\begin{equation*}
B_{\mathbf{Y}} \subseteq S B_{\mathbf{X}} \quad \text { and } \quad \frac{\mathrm{iq}\left(B_{\mathbf{Y}}\right)}{\sqrt{m_{1}+\ldots+m_{n}}}\left(\frac{\operatorname{vol}_{m_{1}+\ldots+m_{n}}\left(B_{\mathbf{X}}\right)}{\mathrm{vol}_{m_{1}+\ldots+m_{n}}\left(B_{\mathbf{Y}}\right)}\right)^{\frac{1}{m_{1}+\ldots+m_{n}}} \lesssim \alpha . \tag{320}
\end{equation*}
$$

Iffurthermore $S_{1}, \ldots, S_{n}$ are all identity mappings (of the respective dimensions), then $S$ can be taken to be the identity mapping provided the the following two conditions hold:

$$
\begin{equation*}
\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{E}}\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{E}^{*}} \lesssim n \tag{321}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(\prod_{k=1}^{n} m_{k}^{m_{k}} \operatorname{vol}_{m_{k}}\left(B_{\mathbf{X}_{k}}\right)\right)^{\frac{1}{m_{1}+\ldots+m_{n}}} \lesssim \frac{m_{1}+\ldots+m_{n}}{n} \min _{k \in\{1, \ldots, n\}} \operatorname{vol}_{m_{k}}\left(B \mathbf{X}_{k}\right)^{\frac{1}{m_{k}}} . \tag{322}
\end{equation*}
$$

Note that (322) is satisfied in particular if $m_{i}=m_{j}$ and $\operatorname{vol}_{m_{i}}\left(B_{\mathbf{X}_{i}}\right)^{\frac{1}{m_{i}}}=\operatorname{vol}_{m_{i}}\left(B_{\mathbf{X}_{j}}\right)^{\frac{1}{m_{j}}}$ for every $i, j \in\{1, \ldots, n\}$.
Prior to proving (150) we will make some basic observations. Firstly, 318 indeed defines a norm because it is well-known that the requirement that $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ is an unconditional normed space is equivalent to (see e.g. LT77, Proposition 1.c.7]) the following "contraction property."

$$
\begin{equation*}
\forall a, x \in \mathbb{R}^{n}, \quad\left\|\left(a_{1} x_{1}, \ldots, a_{n} x_{n}\right)\right\|_{\mathbf{E}} \leqslant\|a\|_{\ell_{\infty}^{n}}\|x\|_{\mathbf{E}} \tag{323}
\end{equation*}
$$

Thus, $\|x\|_{\mathbf{E}} \leqslant\|y\|_{\mathbf{E}}$ if $x, y \in \mathbb{R}^{n}$ satisfy $\left|x_{i}\right| \leqslant\left|y_{i}\right|$ for every $i \in\{1, \ldots, n\}$, so the triangle inequality for (318) follows from applying the triangle inequalities entry-wise for each of the norms $\left\{\|\cdot\|_{\mathbf{x}_{i}}\right\}_{i=1}^{n}$, using this monotonicity property, and then applying the triangle inequality for $\|\cdot\|_{\mathbf{E}}$.

It is well-known that condition (321) holds (as an equality) when $\mathbf{E}$ is a symmetric normed space (see e.g. [LT79, Proposition 3.a.6]). More generally, condition (321] holds (also as an equality) in the setting of the following simple averaging lemma, which shows in particular that Lemma 150 implies Lemma 52 .
Lemma 151. Suppose that $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is a normed space such that for every $j, k \in\{1, \ldots, n\}$ there exists a permutation $\pi \in S_{n}$ with $\pi(j)=k$ such that $\left\|\sum_{i=1}^{n} a_{\pi(i)} e_{i}\right\|_{\mathbf{x}}=\left\|\sum_{i=1}^{n} a_{i} e_{i}\right\|_{\mathbf{x}}$ for every $a_{1}, \ldots, a_{n} \in \mathbb{R}$. Then,

$$
\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{X}}\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{X}^{*}}=n
$$

Proof. Denote $\mathfrak{S}(\mathbf{X})=\left\{\pi \in S_{n}: T_{\pi} \in \operatorname{Isom}(\mathbf{X})\right\}$, where $T_{\pi} \in G L_{n}(\mathbb{R})$ was defined in Example 39 for each $\pi \in S_{n}$. Then, $\mathfrak{S}(\mathbf{X})$ is a subgroup of $S_{n}$ that we are assuming acts transitively on $\{1, \ldots, n\}$. Consequently,

$$
\begin{equation*}
\forall i, j \in\{1, \ldots, n\}, \quad|\{\pi \in \mathfrak{S}(\mathbf{X}): \pi(i)=j\}|=\frac{|\mathfrak{S}(\mathbf{X})|}{n} \tag{324}
\end{equation*}
$$

For every $a_{1}, \ldots, a_{n} \in \mathbb{R}$ we have

$$
\frac{1}{|\mathfrak{S}(\mathbf{X})|} \sum_{\pi \in \mathfrak{S}(\mathbf{X})} \sum_{i=1}^{n} a_{\pi(i)} e_{i}=\sum_{i=1}^{n}\left(\sum_{j=1}^{n} \frac{|\{\pi \in \mathfrak{S}(\mathbf{X}): \pi(i)=j\}|}{|\mathfrak{S}(\mathbf{X})|} a_{j}\right) e_{i} \stackrel{[324}{-} \frac{\sum_{j=1}^{n} a_{j}}{n} \sum_{i=1}^{n} e_{i}
$$

Hence,

$$
\begin{aligned}
\left|\left\langle\sum_{j=1}^{n} e_{j}, \sum_{j=1}^{n} a_{j} e_{j}\right\rangle\right|=\left|\sum_{j=1}^{n} a_{j}\right| & =\frac{n\left\|\frac{1}{|\mathfrak{S}(\mathbf{X})|} \sum_{\pi \in \mathfrak{S}(\mathbf{X})} \sum_{i=1}^{n} a_{\pi(i)} e_{i}\right\|_{\mathbf{X}}}{\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{X}}} \\
& \leqslant \frac{\frac{n}{|\mathfrak{S}(\mathbf{X})|} \sum_{\pi \in \mathfrak{S}(\mathbf{X})}\left\|\sum_{i=1}^{n} a_{\pi(i)} e_{i}\right\|_{\mathbf{X}}}{\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{X}}}=\frac{n\left\|\sum_{i=1}^{n} a_{i} e_{i}\right\|_{\mathbf{X}}}{\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{X}}}
\end{aligned}
$$

where the penultimate step uses convexity and the final step uses the assumption that $T_{\pi}$ is an isometry of $\mathbf{X}$ for every $\pi \in \mathfrak{S}(\mathbf{X})$. Since this holds for every $a_{1}, \ldots, a_{n} \in \mathbb{R}$, we have $\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{X}^{*}} \leqslant n /\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{X}}$. The reverse inequality holds for any normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ because $\left\langle\sum_{i=1}^{n} e_{i}, \sum_{i=1}^{n} e_{i}\right\rangle=n$.

By combining Lemma 150 and Lemma 151 we obtain the following corollary that establishes Conjecture 48 for the iteratively nested $\ell_{p}$ spaces of Lemma 149 , provided it holds for the initial space $\mathbf{X}$.
Corollary 152. Fix $\left\{n_{k}\right\}_{k=0}^{\infty} \subseteq \mathbb{N}$ and $\left\{p_{k}\right\}_{k=1}^{\infty} \subseteq[1, \infty]$. Let $\mathbf{X}=\left(\mathbb{R}^{n_{0}},\|\cdot\|_{\mathbf{X}}\right)$ be a normed space and define

$$
\begin{equation*}
\forall k \in \mathbb{N}, \quad \mathbf{X}_{k+1}=\ell_{p_{k}}^{n_{k}}\left(\mathbf{X}_{k}\right), \quad \text { where } \quad \mathbf{X}_{0}=\mathbf{X} \tag{325}
\end{equation*}
$$

Suppose that $\alpha>0$ and there exists a normed space $\mathbf{Y}=\left(\mathbb{R}^{n_{0}},\|\cdot\|_{\mathbf{Y}}\right)$ with $B_{\mathbf{Y}} \subseteq B_{\mathbf{X}}$ and that satisfies

$$
\begin{equation*}
\frac{\mathrm{iq}\left(B_{\mathbf{Y}}\right)}{\sqrt{n_{0}}}\left(\frac{\operatorname{vol}_{n_{0}}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{n_{0}}\left(B_{\mathbf{Y}}\right)}\right)^{\frac{1}{n_{0}}} \leqslant \alpha . \tag{326}
\end{equation*}
$$

Then, for every $m \in \mathbb{N}$ there is a normed space $\mathbf{Y}_{m}=\left(\mathbb{R}^{n_{0} \cdots n_{m}},\|\cdot\|_{\mathbf{Y}_{m}}\right)$ with $B_{\mathbf{Y}_{m}} \subseteq B_{\mathbf{X}_{m}}$ and

$$
\frac{\operatorname{iq}\left(B_{\mathbf{Y}_{m}}\right)}{\sqrt{n_{0} \cdots n_{m}}}\left(\frac{\operatorname{vol}_{n_{0} \cdots n_{m}}\left(B_{\mathbf{X}_{m}}\right)}{\operatorname{vol}_{n_{0} \cdots n_{m}}\left(B_{\mathbf{Y}_{m}}\right)}\right)^{\frac{1}{n_{0} \cdots n_{m}}} \lesssim \alpha
$$

To see why Corollary 152 indeed follows from Lemma 150 and Lemma 151 observe that if we start with $\mathbf{E}_{0}=\mathbb{R}$ and define inductively $\mathbf{E}_{k+1}=\ell_{p_{k}}^{n_{k}}\left(\mathbf{E}_{k}\right)$, then for each $m \in \mathbb{N}$ the space $\mathbf{E}_{m}$ is unconditional and satisfies the assumptions of Lemma 151. The space $\mathbf{Y}_{m}$ of Corollary 152 is the same space that is defined in Lemma 150 if we take $\mathbf{E}=\mathbf{E}_{m}$, and also $\mathbf{X}_{1}=\ldots=\mathbf{X}_{m}=\mathbf{X}$, which ensures that (322) holds.

Proof of Lemma 150 . Denote

$$
\begin{equation*}
M \stackrel{\text { def }}{=} \sum_{k=1}^{n} m_{k}=\operatorname{dim}(\mathbf{X}) \quad \text { and } \quad \forall k \in\{1, \ldots, n\}, \quad \rho_{k} \stackrel{\text { def }}{=} \operatorname{vol}_{m_{k}}\left(B_{\mathbf{X}_{k}}\right)^{\frac{1}{m_{k}}} . \tag{327}
\end{equation*}
$$

Fix positive numbers $c, C_{1}, \ldots, C_{n}, \gamma_{1}, \ldots, \gamma_{n}, w_{1}, \ldots w_{n}, w_{1}^{*}, \ldots, w_{n}^{*}, \beta_{1}, \ldots, \beta_{n}>0$ that satisfy the following conditions (their values will be specified later). Firstly, we require that

$$
\begin{equation*}
\left\|\sum_{i=1}^{n} w_{i} e_{i}\right\|_{\mathbf{E}}=\left\|\sum_{i=1}^{n} w_{i}^{*} e_{i}\right\|_{\mathbf{E}^{\mathbf{t}}}=1 \tag{328}
\end{equation*}
$$

Secondly, we require that

$$
\begin{equation*}
\forall k \in\{1, \ldots, n\}, \quad w_{k} w_{k}^{*} \geqslant \frac{m_{k}}{\gamma_{k} M} . \tag{329}
\end{equation*}
$$

Finally, we require that

$$
\begin{equation*}
\forall k \in\{1, \ldots, n\}, \quad \frac{1}{c w_{k} \rho_{k}} \leqslant \beta_{k} \leqslant \frac{C_{k}}{w_{k} \rho_{k}}, \tag{330}
\end{equation*}
$$

Denote

$$
\begin{equation*}
D \stackrel{\text { def }}{=}\left(\prod_{k=1}^{n} \beta_{k}^{m_{k}}\right)^{\frac{1}{M}} . \tag{331}
\end{equation*}
$$

Consider the block diagonal linear operator $S: \mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}} \rightarrow \mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}}$ that is given by

$$
\begin{equation*}
\forall x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}}, \quad S x \stackrel{\text { def }}{=} \frac{1}{D}\left(\beta_{1} S_{1} x_{1}, \ldots, \beta_{n} S_{n} x_{n}\right) . \tag{332}
\end{equation*}
$$

The normalization by $D$ in (332) ensures that $S \in \operatorname{SL}\left(\mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}}\right)$.
Since $\sum_{k=1}^{n} w_{k}^{*} e_{k}$ is a unit functional in $\mathbf{E}^{*}$, for every $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}}$ we have

This shows that

Using Lemma 148 , we therefore have

$$
\begin{align*}
\operatorname{vol}_{M}(B \mathbf{X})^{\frac{1}{M}} & \leqslant \frac{M}{D} \operatorname{vol}_{M}\left(B\left(\frac{\gamma_{1} w_{1} \beta_{1}}{m_{1}} S_{1} \mathbf{x}_{1}\right) \oplus_{\oplus_{1} \ldots \oplus_{1}}\left(\frac{\gamma_{n} w_{n} \beta_{n}}{m_{n}} S_{n} \mathbf{x}_{n}\right)\right)^{\frac{1}{M}} \stackrel{314}{=} \frac{1}{D}\left(\frac{M^{M}}{M!} \prod_{k=1}^{n} m_{k}!\left(\frac{\gamma_{k} w_{k} \beta_{k} \rho_{k}}{m_{k}}\right)^{m_{k}}\right)^{\frac{1}{M}}  \tag{334}\\
& \stackrel{\sqrt{330}}{\leqslant} \frac{1}{D}\left(\frac{M^{M}}{M!} \prod_{k=1}^{n} \frac{m_{k}!}{m_{k}^{m_{k}}}\left(\gamma_{k} C_{k}\right)^{m_{k}}\right)^{\frac{1}{M}} \leqslant \frac{e}{D}\left(\prod_{k=1}^{n}\left(\gamma_{k} C_{k}\right)^{m_{k}}\right)^{\frac{1}{M}} .
\end{align*}
$$

Next, for every $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}}$ we have

$$
\begin{aligned}
\left\|S^{-1} x\right\|_{\mathbf{X}} \stackrel{\sqrt[318]{ }}{=} \frac{\sqrt{332}}{} D & \left\|\sum_{k=1}^{n} \frac{\left\|S_{k}^{-1} x_{k}\right\|_{\mathbf{x}_{k}}}{\beta_{k}} e_{k}\right\|_{\mathbf{E}} \\
& \stackrel{\boxed{323}}{\leqslant} D\left(\max _{k \in\{1, \ldots, n\}} \frac{\left\|S_{k}^{-1} x_{k}\right\|_{\mathbf{x}_{k}}}{w_{k} \beta_{k}}\right)\left\|\sum_{k=1}^{n} w_{k} e_{k}\right\|_{\mathbf{E}} \frac{\sqrt{328}}{-} D \max _{k \in\{1, \ldots, n\}} \frac{\left\|S_{k}^{-1} x_{k}\right\|_{\mathbf{x}_{k}}}{w_{k} \beta_{k}} .
\end{aligned}
$$

This establishes the following inclusion.

$$
\begin{equation*}
S B_{\mathbf{X}} \supseteq \frac{1}{D} \prod_{k=1}^{n} w_{k} \beta_{k} S_{k} B_{\mathbf{X}_{k}} \stackrel{\text { def }}{=} \Omega \tag{335}
\end{equation*}
$$

Thanks to (62), the assumption (319) of Lemma 150 implies that

$$
\begin{equation*}
\forall k \in\{1, \ldots, n\}, \quad \lambda\left(S_{k} B_{\mathbf{X}_{k}}\right) \rho_{k}^{2} \stackrel{[327}{-} \lambda\left(S_{k} B_{\mathbf{X}_{k}}\right) \operatorname{vol}_{m_{k}}\left(B_{\mathbf{X}_{k}}\right)^{\frac{2}{m_{k}}} \lesssim \alpha^{2} m_{k} \tag{336}
\end{equation*}
$$

For each $k \in\{1, \ldots, n\}$ take $f_{k}: S_{k} B_{\mathbf{X}_{k}} \rightarrow \mathbb{R}$ that is smooth on the interior of $S_{k} B_{\mathbf{X}_{k}}$, vanishes on $\partial S_{k} B_{\mathbf{X}_{k}}$, and satisfies $\Delta f_{k}=-\lambda\left(S_{k} B_{\mathbf{X}_{k}}\right) f_{k}$ on the interior of $S_{k} B_{\mathbf{X}_{k}}$. Define $f: \Omega \rightarrow \mathbb{R}$ by

$$
\forall x=\left(x_{1}, \ldots, x_{n}\right) \in \Omega=\frac{1}{D} \prod_{k=1}^{n} w_{k} \beta_{k} S_{k} B_{\mathbf{X}_{k}}, \quad f(x) \stackrel{\text { def }}{=} \prod_{k=1}^{n} f_{k}\left(\frac{D}{w_{k} \beta_{k}} x_{k}\right)
$$

Thus $f \equiv 0$ on the boundary of $\Omega$ and on the interior of $\Omega$ it is smooth and satisfies

$$
\begin{equation*}
\Delta f=-D^{2}\left(\sum_{k=1}^{n} \frac{\lambda\left(S_{k} B \mathbf{X}_{k}\right)}{\left(w_{k} \beta_{k}\right)^{2}}\right) f \tag{337}
\end{equation*}
$$

Hence,

$$
\begin{equation*}
\lambda(S \mathbf{X})=\lambda\left(S B_{\mathbf{X}}\right) \stackrel{\sqrt{335}}{\leqslant} \lambda(\Omega) \stackrel{\sqrt[337]{\leqslant}}{\leqslant} D^{2}\left(\sum_{k=1}^{n} \frac{\lambda\left(S_{k} B_{\mathbf{X}_{k}}\right)}{\left(w_{k} \beta_{k}\right)^{2}}\right) \stackrel{\sqrt{330}}{\leqslant}(c D)^{2}\left(\sum_{k=1}^{n} \lambda\left(S_{k} B_{\mathbf{X}_{k}}\right) \rho_{k}^{2}\right) \stackrel{\sqrt[336]{\lesssim}}{\lesssim}(c \alpha D)^{2} M \tag{338}
\end{equation*}
$$

By combining (334) and (338) we see that

$$
\lambda(S \mathbf{X}) \operatorname{vol}_{M}\left(B_{\mathbf{X}}\right)^{\frac{2}{M}} \lesssim c^{2}\left(\prod_{k=1}^{n}\left(\gamma_{k} C_{k}\right)^{m_{k}}\right)^{\frac{2}{M}} \alpha^{2} M
$$

Another application of (62) now shows that the desired conclusion (320) holds with $\mathbf{Y}=\mathrm{ChSX}$ (recall the definition of Cheeger space in Section 1.6.1) provided

$$
\begin{equation*}
c\left(\prod_{k=1}^{n}\left(\gamma_{k} C_{k}\right)^{m_{k}}\right)^{\frac{1}{M}} \lesssim 1 \tag{339}
\end{equation*}
$$

To get (320), by the Lozanovskiĭ factorization theorem [Loz69] there exist $w_{1}, \ldots w_{n}, w_{1}^{*}, \ldots, w_{n}^{*}>0$ such that 328 holds and also $w_{k} w_{k}^{*}=m_{k} / M$ for every $k \in\{1, \ldots, n\}$. Thus (329) holds (as an equality) if we choose $\gamma_{1}=\ldots=\gamma_{n}=1$. If we take $c=C_{1}=\ldots=C_{n}=1$ and $\beta_{k}=1 /\left(w_{k} \rho_{k}\right)$ for each $k \in\{1, \ldots, n\}$, then both (330) and (339) also hold (as equalities). With these choices, 320) holds.

Suppose that the additional assumptions (321) and 322) hold. Denote $\eta=\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{E}}\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{E}^{*}} / n$. So, $\eta=O(1)$ by (321). Take $w_{1}=\ldots=w_{n}=1 /\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{E}}$ and $w_{1}^{*}=\ldots=w_{n}^{*}=1 /\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{E}^{*}}$, so that (328) holds by design. This choice also ensures that if we take $\gamma_{k}=m_{k} /(\eta M)$ for each $k \in\{1, \ldots, n\}$, then (329) holds (as an equality). Next, choose $C_{k}=\rho_{k}$ for each $k \in\{1, \ldots, n\}$, as well as $\beta_{1}=\ldots=\beta_{n}=\left\|\sum_{i=1}^{n} e_{i}\right\|_{\mathbf{E}}$ and $c=1 / \min _{k \in\{1, \ldots, n\}} \rho_{k}$. This ensures that (330) holds, and also that (339) coincides with the assumption (322), since $\eta=O(1)$. The desired conclusion (320) therefore holds with $S x=\left(S_{1} x_{1}, \ldots, S_{n} x_{n}\right)$ in (332). In particular, if $S_{k}=\operatorname{ld}_{m_{k}}$ for every $k \in\{1, \ldots, n\}$, then we can take $S=\operatorname{ld}_{\mathbb{R}^{m_{1}} \times \ldots \times \mathbb{R}^{m_{n}}}$ in (320).

The following lemma provides a formula for the cone measure of Orlicz spaces. Fix a convex increasing function $\psi:[0, \infty) \rightarrow[0, \infty]$ that satisfies $\psi(0)=0$ and $\lim _{x \rightarrow \infty} \psi(x)=\infty$ (so, if $\lim _{x \rightarrow a^{-}} \psi(x)=\infty$ for some $a \in(0, \infty)$, then we require that $\psi(x)=\infty$ for every $x \geqslant a)$. Henceforth, the associated Orlicz space (see e.g. RR91b]) $\ell_{\psi}^{n}=\left(\mathbb{R}^{n},\|\cdot\|_{\ell_{\psi}^{n}}\right)$ will always be endowed with the Luxemburg norm that is given by

$$
\begin{equation*}
\forall x \in \mathbb{R}^{n}, \quad\|x\|_{\ell_{\psi}^{n}}=\inf \left\{s>0: \sum_{i=1}^{n} \psi\left(\frac{\left|x_{i}\right|}{s}\right) \leqslant 1\right\} \tag{340}
\end{equation*}
$$

Lemma 153. Suppose that $\psi:[0, \infty) \rightarrow[0, \infty]$ is convex, increasing, continuously differentiable on the set $\{x \in(0, \infty): \psi(x)<\infty\}$, and satisfies $\psi(0)=0$ and $\lim _{x \rightarrow \infty} \psi(x)=\infty$. Then, for every $g \in L_{1}\left(\kappa_{\ell_{\psi}^{n}}^{n}\right.$ we have

$$
\begin{align*}
& \frac{n!}{2^{n}} \operatorname{vol}_{n}\left(B_{\ell_{\psi}^{n}}\right) \int_{\partial B_{\ell_{\psi}^{n}}} g(\theta) \mathrm{d} \kappa_{\ell_{\psi}^{n}}^{n}(\theta) \\
& \quad=\int_{\partial B_{\ell_{1}^{n}}} g\left(\psi^{-1}\left(\left|\tau_{i}\right|\right) \operatorname{sign}\left(\tau_{1}\right), \ldots, \psi^{-1}\left(\left|\tau_{n}\right|\right) \operatorname{sign}\left(\tau_{n}\right)\right) \frac{\sum_{i=1}^{n} \psi^{-1}\left(\left|\tau_{i}\right|\right) \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)}{\prod_{i=1}^{n} \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)} \mathrm{d} \kappa_{\ell_{1}^{n}}(\tau) . \tag{341}
\end{align*}
$$

For example, when $\psi(t)=t^{p}$ for some $p \geqslant 1$ and every $t \geqslant 0$, in which case $\ell_{\psi}^{n}=\ell_{p}^{n}$, Lemma 153 gives

$$
\int_{\partial B_{\ell_{p}^{n}}} g \mathrm{~d} \kappa_{\ell_{\psi}^{n}}=\frac{\Gamma\left(1+\frac{n}{p}\right)}{n!\Gamma\left(1+\frac{1}{p}\right)^{n}} \int_{\partial B_{\ell_{1}^{n}}} \frac{g \circ M_{1 \rightarrow p}^{n}(\tau)}{\left|\tau_{1} \cdots \tau_{n}\right|^{1-\frac{1}{p}}} \mathrm{~d} \kappa_{\ell_{1}^{n}}(\tau),
$$

where $M_{1 \rightarrow p}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is the Mazur map [Maz29] from $\ell_{1}^{n}$ to $\ell_{p}^{n}$, i.e.,

$$
\forall x \in \mathbb{R}^{n}, \quad M_{1 \rightarrow p}^{n}\left(x_{1}, \ldots, x_{n}\right)=\left(\left|x_{1}\right|^{\frac{1}{p}} \operatorname{sign}\left(x_{1}\right), \ldots,\left|x_{n}\right|^{\frac{1}{p}} \operatorname{sign}\left(x_{n}\right)\right) .
$$

As another special case of Lemma 153 , consider the following family of Orlicz spaces $\Omega_{\beta}^{n}=\left(\mathbb{R}^{n},\|\cdot\|_{\Omega_{\beta}^{n}}\right)$ :

$$
\forall \beta>0, \quad \Omega_{\beta}^{n} \stackrel{\text { def }}{=} \ell_{\psi_{\beta}}^{n} \quad \text { where } \quad \forall t \geqslant 0, \quad \psi_{\beta}(t) \stackrel{\text { def }}{=} \begin{cases}\frac{1}{\beta} \log \left(\frac{1}{1-t}\right) & \text { if } 0 \leqslant t<1,  \tag{342}\\ \infty & \text { if } t \geqslant 1 .\end{cases}
$$

Observe that by considering the case $g \equiv 1$ of (341) we obtain the following identity.

$$
\begin{equation*}
\int_{\partial B_{\ell_{\psi}}} g \mathrm{~d} \kappa_{\ell_{\psi}^{n}}=\frac{\int_{\partial B_{\ell_{1}^{n}}} g\left(\psi^{-1}\left(\left|\tau_{i}\right|\right) \operatorname{sign}\left(\tau_{1}\right), \ldots, \psi^{-1}\left(\left|\tau_{n}\right|\right) \operatorname{sign}\left(\tau_{n}\right)\right) \frac{\sum_{i=1}^{n} \psi^{-1}\left(\left|\tau_{i}\right|\right) \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)}{\prod_{i=1}^{n} \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)} \mathrm{d} \kappa_{\ell_{1}^{n}}(\tau)}{\int_{\partial B_{\ell_{1}^{n}}} \frac{\sum_{i=1}^{n} \psi^{-1}\left(\left|\tau_{i}\right| \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)\right.}{\prod_{i=1}^{n} \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)} \mathrm{d} \kappa_{\ell_{1}^{n}}(\tau)} . \tag{343}
\end{equation*}
$$

When $\psi=\psi_{\beta}$ for some $\beta>0$ (we will eventually need to work with $\beta=n$ ), for every $\tau \in \partial B_{\ell_{1}^{n}}$ we have

$$
\begin{equation*}
\frac{\sum_{i=1}^{n} \psi_{\beta}^{-1}\left(\left|\tau_{i}\right|\right) \psi_{\beta}^{\prime}\left(\psi_{\beta}^{-1}\left(\left|\tau_{i}\right|\right)\right)}{\prod_{i=1}^{n} \psi_{\beta}^{\prime}\left(\psi_{\beta}^{-1}\left(\left|\tau_{i}\right|\right)\right)}=\frac{\sum_{i=1}^{n}\left(1-e^{-\beta\left|\tau_{i}\right|}\right) \frac{e^{\beta\left|\tau_{i}\right|}}{\beta}}{\prod_{i=1}^{n} \frac{e^{\beta\left|\tau_{i}\right|}}{\beta}}=\frac{\beta^{n-1} \sum_{i=1}^{n}\left(e^{\beta\left|\tau_{i}\right|}-1\right)}{e^{\beta\|\tau\| \ell_{1}^{n}}}=\frac{\beta^{n-1}}{e^{\beta}} \sum_{i=1}^{n}\left(e^{\beta\left|\tau_{i}\right|}-1\right) . \tag{344}
\end{equation*}
$$

Consequently, (343) gives the following identity, which we will need later.

$$
\begin{equation*}
\int_{\partial{\Omega_{\Omega}^{n}}_{\beta}} g \mathrm{~d} \kappa_{\Omega_{\beta}^{n}}=\frac{\int_{\partial B_{\ell_{1}^{n}}} g\left(\left(e^{\beta\left|\tau_{1}\right|}-1\right) \operatorname{sign}\left(\tau_{1}\right), \ldots,\left(e^{\beta\left|\tau_{n}\right|}-1\right) \operatorname{sign}\left(\tau_{n}\right)\right) \sum_{i=1}^{n}\left(e^{\beta\left|\tau_{i}\right|}-1\right) \mathrm{d} \kappa_{\ell_{1}^{n}}(\tau)}{\int_{\partial B_{\ell_{1}^{n}}} \sum_{i=1}^{n}\left(e^{\beta\left|\tau_{i}\right|}-1\right) \mathrm{d} \kappa_{\ell_{1}^{n}}(\tau)} . \tag{345}
\end{equation*}
$$

Proof of Lemma 153 . For each $i \in\{1, \ldots, n\}$ define $f_{i}: \mathbb{R}^{n} \rightarrow \mathbb{R}$ by setting $f_{i}(0)=0$ and

$$
\forall y \in \mathbb{R}^{n} \backslash\{0\}, \quad f_{i}(y)=\|y\|_{\ell_{1}^{n}} \psi^{-1}\left(\frac{\left|y_{i}\right|}{\|y\|_{\ell}^{n}}\right) \operatorname{sign}\left(y_{i}\right)
$$

Consider $f=\left(f_{1}, \ldots, f_{n}\right): \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$. Then, $\|f(y)\|_{\ell_{\psi}^{n}}=\|y\|_{\ell_{1}^{n}}$ for every $y \in \mathbb{R}^{n}$. Hence, $f\left(B_{\ell_{1}^{n}}\right)=B \ell_{\psi}^{n}$. Now,

$$
\begin{aligned}
& \int_{\partial B_{\ell_{\psi}^{n}}} g(\theta) \mathrm{d} \kappa_{\ell_{\psi}^{n}}(\theta) \stackrel{\sqrt[333]{=}}{=} \frac{1}{\operatorname{vol}_{n}\left(B_{\ell_{\psi}^{n}}\right)} \int_{f\left(B_{\ell_{1}^{n}}\right)} g\left(\frac{1}{\|x\|_{\ell_{\psi}^{n}}} x\right) \mathrm{d} x \\
& \quad=\frac{1}{\operatorname{vol}_{n}\left(B_{\ell_{\psi}^{n}}\right)} \int_{B_{\ell_{1}^{n}}} g\left(\frac{1}{\|f(y)\|_{\ell_{\psi}^{n}}} f(y)\right)\left|\operatorname{det} f^{\prime}(y)\right| \mathrm{d} y \stackrel{(313}{=} \frac{\operatorname{vol}_{n}\left(B_{\ell_{1}^{n}}\right)}{\operatorname{vol}_{n}\left(B_{\ell_{\psi}}^{n}\right.} \int_{\partial B_{\ell_{1}^{n}}} g(f(\tau))\left|\operatorname{det} f^{\prime}(\tau)\right| \mathrm{d} \kappa_{\ell_{1}^{n}}(\tau),
\end{aligned}
$$

where in the final step we used the fact $f$ is positively homogeneous of order 1 , and hence its derivative is homogeneous of order 0 almost everywhere ( $f$ is continuously differentiable on $\left\{y \in \mathbb{R}^{n} ; y_{1}, \ldots, y_{n} \neq 0\right\}$ ). Since the volume of the unit ball of $\ell_{1}^{n}$ equals $2^{n} / n!$, it remains to check that the Jacobian of $f$ satisfies

$$
\operatorname{det} f^{\prime}(\tau)=\frac{\sum_{i=1}^{n} \psi^{-1}\left(\left|\tau_{i}\right|\right) \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)}{\prod_{i=1}^{n} \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)},
$$

for every $\tau \in \partial B_{\ell_{1}^{n}}$ with $\tau_{1}, \ldots, \tau_{n} \neq 0$. This is so because for every such $\tau$ and $i, j \in\{1, \ldots, n\}$ we have

$$
\partial_{j} f_{i}(\tau)=\frac{\delta_{i j}-\tau_{i} \operatorname{sign}\left(\tau_{j}\right)}{\psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)}+\psi^{-1}\left(\left|\tau_{i}\right|\right) \operatorname{sign}\left(\tau_{i}\right) \operatorname{sign}\left(\tau_{j}\right) .
$$

Hence, $f^{\prime}(\tau)=A(\tau)+u(\tau) \otimes v(\tau)$, where $A(\tau) \in \mathrm{M}_{n}(\mathbb{R})$ is the diagonal matrix $\operatorname{Diag}\left(\left(1 / \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)\right)_{i=1}^{n}\right)$ and $u(\tau)=\left(\psi^{-1}\left(\left|\tau_{i}\right|\right) \operatorname{sign}\left(\tau_{i}\right)-\tau_{i} / \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)\right)_{i=1}^{n}, \nu(\tau)=\left(\operatorname{sign}\left(\tau_{i}\right)\right)_{i=1}^{n} \in \mathbb{R}^{n}$. By the textbook formula for the determinant of a rank-1 perturbation of an invertible matrix (e.g. Mey00. Section 6.2]), it follows that

$$
\begin{aligned}
& \operatorname{det} f^{\prime}(\tau)=\left(1+\left\langle A(\tau)^{-1} u(\tau), v(\tau)\right\rangle\right) \operatorname{det} A(\tau) \\
& =\frac{1+\sum_{i=1}^{n} \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)\left(\psi^{-1}\left(\left|\tau_{i}\right|\right) \operatorname{sign}\left(\tau_{i}\right)-\frac{\tau_{i}}{\psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)}\right) \operatorname{sign}\left(\tau_{i}\right)}{\prod_{i=1}^{n} \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)}=\frac{\sum_{i=1}^{n} \psi^{-1}\left(\left|\tau_{i}\right|\right) \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)}{\prod_{i=1}^{n} \psi^{\prime}\left(\psi^{-1}\left(\left|\tau_{i}\right|\right)\right)} .
\end{aligned}
$$

Another description of $\kappa_{\mathbf{x}}$ is the fact (see e.g. [NR03, Lemma 1]) that the Radon-Nikodým derivative of the ( $n-1$ )-dimensional Hausdorff (non-normalized surface area) measure on $\partial B_{\mathbf{X}}$ with respect to the (non-normalized cone) measure $\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)_{\mathbf{X}}$ is equal at almost every $x \in \partial B_{\mathbf{X}}$ to $n$ times the Euclidean length of the gradient at $x$ of the function $u \mapsto\|u\|_{\mathbf{X}}$. In other words, for any $g \in L_{1}\left(\partial B_{\mathbf{X}}\right)$,

$$
\begin{equation*}
\int_{\partial B_{\mathbf{X}}} g(x) \mathrm{d} x=n \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) \int_{\partial B_{\mathbf{X}}} g(x)\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}} \mathrm{~d} \kappa_{\mathbf{X}}(x) . \tag{346}
\end{equation*}
$$

The special case $g \equiv 1$ of (346) gives the following identity.

$$
\begin{equation*}
\frac{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}=n \int_{\partial B_{\mathbf{X}}}\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}} \mathrm{~d} \kappa_{\mathbf{X}}(x)=f_{B_{\mathbf{X}}} \frac{\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}}}{\|x\|_{\mathbf{X}}^{n-1}} \mathrm{~d} x, \tag{347}
\end{equation*}
$$

where the second equality in (347) is an application of (313) because it is straightforward to check that $\|\nabla\| \cdot\left\|_{\mathbf{X}}(r x)\right\|_{\ell_{2}^{n}}=\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}}$ for any $r>0$ and $x \in \mathbb{R}^{n}$ at which the norm $\|\cdot\|_{\mathbf{X}}$ is smooth.
Remark 154. By applying Cauchy-Schwarz to the first equality in (347), we see that

$$
\begin{equation*}
\frac{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} \leqslant n\left(\int_{\partial B_{\mathbf{X}}}\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}}^{2} \mathrm{~d} \kappa_{\mathbf{X}}(x)\right)^{\frac{1}{2}}=\left(\frac{n}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} \int_{\partial B_{\mathbf{X}}}\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}} \mathrm{~d} x\right)^{\frac{1}{2}}, \tag{348}
\end{equation*}
$$

where the final step of (348) is an applications of (346) with $g(x)=\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}}$. If $\|\cdot\|_{\mathbf{X}}$ is twice continuously differentiable on $\mathbb{R}^{n} \backslash\{0\}$ and $\varphi: \mathbb{R} \rightarrow[0, \infty)$ is twice continuously differentiable with $\varphi^{\prime}(1)>0$ and $\varphi^{\prime \prime}(0)=0$, then because for every $x \in \partial B_{\mathbf{X}}$ the vector $\nabla\|\cdot\|_{\mathbf{X}}(x) /\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}}$ is the unit outer normal to $\partial B_{\mathbf{X}}$ at $x$, by the divergence theorem we have

$$
\begin{aligned}
& \int_{\partial B_{\mathbf{X}}} \Delta\left(\varphi \circ\|\cdot\|_{\mathbf{X}}\right)(x) \mathrm{d} x=\int_{\partial B_{\mathbf{X}}} \operatorname{div} \nabla\left(\varphi \circ\|\cdot\|_{\mathbf{X}}\right)(x) \mathrm{d} x=\int_{\partial B_{\mathbf{X}}} \frac{\left\langle\nabla\left(\varphi \circ\|\cdot\|_{\mathbf{X}}\right)(x), \nabla\|\cdot\|_{\mathbf{X}}(x)\right\rangle}{\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}}} \mathrm{~d} x \\
&=\int_{\partial B_{\mathbf{X}}} \varphi^{\prime}\left(\|x\|_{\mathbf{X}}\right)\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}} \mathrm{~d} x=\varphi^{\prime}(1) \int_{\partial B_{\mathbf{X}}}\|\nabla\| \cdot\left\|_{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}} \mathrm{~d} x .
\end{aligned}
$$

A substitution of this identity into (348) give the following bound.

$$
\begin{equation*}
\frac{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} \leqslant \frac{\sqrt{n}}{\sqrt{\varphi^{\prime}(1)}}\left(f_{\partial B_{\mathbf{X}}} \Delta\left(\varphi \circ\|\cdot\|_{\mathbf{X}}\right)(x) \mathrm{d} x\right)^{\frac{1}{2}} . \tag{349}
\end{equation*}
$$

In particular, for every $p>2$ we have

$$
\begin{equation*}
\frac{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} \leqslant \sqrt{\frac{n}{p}}\left(f_{B_{\mathbf{X}}} \Delta\left(\|\cdot\|_{\mathbf{X}}^{p}\right)(x) \mathrm{d} x\right)^{\frac{1}{2}} \tag{350}
\end{equation*}
$$

It is worthwhile to record (349) separately because this estimate is sometimes convenient for getting good bounds on $\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right)$. In particular, by using (350) when $\mathbf{X}$ is an $\ell_{p}$ direct sum one can obtain an alternative derivation of some of the ensuing estimates. Another noteworthy consequence of (348) is when there is a transitive subgroup of permutations $G \leqslant S_{n}$ such that $\left\|\left(x_{\pi(1)}, \ldots, x_{\pi(n)}\right)\right\|_{\mathbf{X}}=\|x\|_{\mathbf{X}}$ for all $x \in \mathbb{R}^{n}$ and $\pi \in G$. Under this further symmetry assumption, the first inequality of (348) becomes

$$
\frac{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)} \leqslant n^{\frac{3}{2}}\left(\int_{\partial B_{\mathbf{X}}}\left(\frac{\partial\|\cdot\|_{\mathbf{X}}}{\partial x_{1}}(x)\right)^{2} \mathrm{~d} \kappa_{\mathbf{X}}(x)\right)^{\frac{1}{2}}
$$

The following lemma provides a probabilistic interpretation of the cone measure which generalizes the treatment of the special case $\mathbf{X}=\ell_{p}^{n}$ by Schechtman-Zinn [SZ90] and Rachev-Rüschendorf [RR91a].

Lemma 155 (probabilistic representation of cone measure). Fix $n \in \mathbb{N}$ and let $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ be a normed space. Suppose that $\varphi:[0, \infty) \rightarrow[0, \infty)$ is a continuous function such that $\varphi(0)=0, \varphi(t)>0$ when $t>0$ and $\int_{0}^{\infty} r^{n-1} \varphi(r) \mathrm{d} r<\infty$. Let $\vee$ be a random vector in $\mathbb{R}^{n}$ whose density at each $x \in \mathbb{R}^{n}$ is equal to

$$
\begin{equation*}
\frac{1}{n \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) \int_{0}^{\infty} r^{n-1} \varphi(r) \mathrm{d} r} \varphi\left(\|x\|_{\mathbf{X}}\right) \tag{351}
\end{equation*}
$$

where we note that (351) in indeed a probability density by (313). Then, the density of $\|\mathrm{V}\|_{\mathrm{X}}$ at $s \in[0, \infty$ ) is equal to $s^{n-1} \varphi(s) / \int_{0}^{\infty} r^{n-1} \varphi(r) \mathrm{d} r$. Moreover, the following two assertions hold:

- $\mathrm{V} /\|\mathrm{V}\|_{\mathbf{X}}$ is distributed according to the cone measure $\kappa_{\mathbf{x}}$,
- $\|\mathrm{V}\|_{\mathbf{X}}$ and $\mathrm{V} /\|\mathrm{V}\|_{\mathbf{X}}$ are (stochastically) independent.

Proof. The density of $\|\mathrm{V}\|_{\mathbf{X}}$ at $s \in[0, \infty)$ is equal to

$$
\begin{gathered}
\frac{\mathrm{d}}{\mathrm{~d} s} \operatorname{Prob}\left[\|\mathrm{~V}\|_{\mathbf{X}} \leqslant s\right] \stackrel{\sqrt[351]{=}}{=} \frac{\mathrm{d}}{\mathrm{~d} s}\left(\frac{1}{n \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) \int_{0}^{\infty} r^{n-1} \varphi(r) \mathrm{d} r} \int_{s B_{\mathbf{X}}} \varphi\left(\|x\|_{\mathbf{X}}\right) \mathrm{d} x\right) \\
\stackrel{313}{-3} \frac{\mathrm{~d}}{\mathrm{~d} s}\left(\frac{\int_{0}^{s} t r^{n-1} \varphi(r) \mathrm{d} r}{\int_{0}^{\infty} r^{n-1} \varphi(r) \mathrm{d} r}\right)=\frac{s^{n-1} \varphi(s)}{\int_{0}^{\infty} r^{n-1} \varphi(r) \mathrm{d} r} .
\end{gathered}
$$

The rest of Lemma 155 is equivalent to showing that for every measurable $A \subseteq \partial B_{\mathbf{X}}$ and $\rho>0$,

$$
\operatorname{Prob}\left[\left.\frac{\mathrm{V}}{\|\mathrm{~V}\|_{\mathbf{X}}} \in A \right\rvert\,\|\mathrm{V}\|_{\mathbf{X}}=\rho\right]=\kappa_{\mathbf{X}}(A)
$$

To prove this identity, observe first that for every $a, b \in \mathbb{R}$ with $a<b$ we have

$$
\operatorname{vol}_{n}([a, b] A)=\operatorname{vol}_{n}\left(b\left(([0,1] A) \backslash\left(\frac{a}{b}[0,1] A\right)\right)\right)=\left(b^{n}-a^{n}\right) \operatorname{vol}_{n}([0,1] A) .
$$

Hence, it follows from the definition (312) that

$$
\begin{equation*}
\kappa_{\mathbf{X}}(A)=\frac{\operatorname{vol}_{n}([a, b] A)}{\operatorname{vol}_{n}\left([a, b] \partial B_{\mathbf{X}}\right)} . \tag{352}
\end{equation*}
$$

Consequently,

$$
\left.\begin{array}{rl}
\operatorname{Prob}\left[\left.\frac{\mathrm{V}}{\|\mathrm{~V}\|_{\mathbf{X}}} \in A \right\rvert\,\|\mathrm{V}\|_{\mathbf{X}}=\right. & \rho]
\end{array}\right)=\lim _{\varepsilon \rightarrow 0} \frac{\operatorname{Prob}\left[\mathrm{~V} \in\|\mathrm{~V}\|_{\mathbf{X}} A \text { and } \rho-\varepsilon \leqslant\|\mathrm{V}\|_{\mathbf{X}} \leqslant \rho+\varepsilon\right]}{\operatorname{Prob}\left[\rho-\varepsilon \leqslant\|\mathrm{V}\|_{\mathbf{X}} \leqslant \rho+\varepsilon\right]}, \quad \begin{array}{ll} 
& \lim _{\varepsilon \rightarrow 0} \frac{\int_{([0, \infty) A) \cap\left([\rho-\varepsilon, \rho+\varepsilon] \partial B_{\mathbf{X}}\right)} \varphi\left(\|x\|_{\mathbf{X}}\right) \mathrm{d} x}{\int_{[\rho-\varepsilon, \rho+\varepsilon] \partial B_{\mathbf{X}}} \varphi\left(\|x\|_{\mathbf{X}}\right) \mathrm{d} x}=\lim _{\varepsilon \rightarrow 0} \frac{\operatorname{vol}_{n}([\rho-\varepsilon, \rho+\varepsilon] A)}{\operatorname{vol}_{n}\left([\rho-\varepsilon, \rho+\varepsilon] \partial B_{\mathbf{X}}\right)}=\kappa_{\mathbf{X}}(A),
\end{array}
$$

where the penultimate step holds as $\varphi$ is continuous at $\rho$ and $\varphi(\rho)>0$, and the final step uses (352).

Lemma 156. Fix $m, n \in \mathbb{N}$ and $p \in(1, \infty)$. Suppose that $\mathbf{X}=\left(\mathbb{R}^{m},\|\cdot\|_{\mathbf{X}}\right)$ is a normed space. Let $\mathrm{R}_{1}, \ldots, \mathrm{R}_{n}$ be i.i.d. random variables taking values in $[0, \infty)$ whose density at each $t \in(0, \infty)$ is equal to

$$
\begin{equation*}
\frac{p}{2(p-1) \Gamma\left(\frac{m}{p}\right)} t^{\frac{m}{2 p-2}-1} e^{-t^{\frac{p}{p p-2}}} . \tag{353}
\end{equation*}
$$

Then,

$$
\begin{equation*}
\frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}(\mathbf{X})}\right)}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}(\mathbf{X})}\right)}=\frac{p \Gamma\left(1+\frac{n m}{p}\right)}{\Gamma\left(1+\frac{n m-1}{p}\right)} \int_{\left(\partial B_{\mathbf{X}}\right)^{n}} \mathbb{E}\left[\left(\sum_{i=1}^{n} \mathrm{R}_{i}\|\nabla\| \cdot\left\|_{\mathbf{X}}\left(x_{i}\right)\right\|_{\ell_{2}^{m}}^{2}\right)^{\frac{1}{2}}\right] \mathrm{d} \kappa_{\mathbf{X}}^{\otimes n}\left(x_{1}, \ldots, x_{n}\right) . \tag{354}
\end{equation*}
$$

Furthermore,

$$
\begin{equation*}
\int_{\partial B_{\ell_{p}^{n}(\mathbf{X})}}\|\nabla\| \cdot\left\|_{\left.\ell_{p}^{n} \mathbf{X}\right)}\right\|_{\ell_{2}^{n}\left(\ell_{2}^{m}\right)}^{2} \mathrm{~d} \kappa_{\ell_{p}^{n}(\mathbf{X})}=\frac{n \Gamma\left(\frac{n m}{p}\right) \Gamma\left(\frac{m+2 p-2}{p}\right)}{\Gamma\left(\frac{m}{p}\right) \Gamma\left(\frac{n m+2 p-2}{p}\right)} \int_{\partial B_{\mathbf{X}}}\|\nabla\| \cdot\left\|_{\mathbf{X}}\right\|_{\ell_{2}^{m}}^{2} \mathrm{~d} \kappa_{\mathbf{X}} . \tag{355}
\end{equation*}
$$

Proof. For almost every $x=\left(x_{1}, \ldots, x_{n}\right) \in \ell_{p}^{n}(\mathbf{X})$ we have

$$
\nabla\|\cdot\|_{\ell_{p}^{n}(\mathbf{X})}(x)=\frac{1}{\|x\|_{\ell_{p}^{n}(\mathbf{X})}^{p-1}}\left(\left\|x_{1}\right\|_{\mathbf{X}}^{p-1} \nabla\|\cdot\|_{\mathbf{X}}\left(x_{1}\right), \ldots,\left\|x_{n}\right\|_{\mathbf{X}}^{p-1} \nabla\|\cdot\|_{\mathbf{X}}\left(x_{n}\right)\right) .
$$

Consequently,

$$
\left.\begin{array}{rl}
\|x\|_{\ell_{p}^{n}(\mathbf{X})}^{p-1}\|\nabla\| \cdot \| \ell_{p}^{n}(\mathbf{X})  \tag{356}\\
\|x\|_{\ell_{p}^{n}(\mathbf{X})}
\end{array}\right) \|_{\ell_{2}^{n}\left(\ell_{2}^{m}\right)}=\left(\sum_{i=1}^{n}\left\|x_{i}\right\|_{\mathbf{X}}^{2 p-2}\|\nabla\| \cdot\left\|\mathbf{X}\left(\frac{x_{i}}{\|x\|_{\ell_{p}^{n}(\mathbf{X})}}\right)\right\|_{\ell_{2}^{m}}^{2}\right)^{\frac{1}{2}}
$$

where we used the straightforward fact that the gradient of any (finite dimensional) norm is homogeneous of order 0 (on its domain of definition, which is almost everywhere).

Let $\mathrm{V}=\left(\mathrm{V}_{1}, \ldots, \mathrm{~V}_{n}\right)$ be a random vector on $\ell_{p}^{n}(\mathbf{X})$ whose density at $x=\left(x_{1}, \ldots, x_{n}\right) \in \ell_{p}^{n}(\mathbf{X})$ is

$$
\begin{equation*}
\frac{1}{\Gamma\left(1+\frac{n m}{p}\right) \operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}(\mathbf{X})}\right)} e^{-\|x\|_{\ell_{p}^{p}(\mathbb{X})}^{p}}=\frac{1}{\Gamma\left(1+\frac{n m}{p}\right) \operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}(\mathbf{X})}\right)} \prod_{i=1}^{n} e^{-\left\|x_{i}\right\|_{\mathrm{X}}^{p}} . \tag{357}
\end{equation*}
$$

By combining Lemma 155 with the first equality in (347), we see that

$$
\begin{equation*}
\frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}(\mathbf{X})}\right)}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}(\mathbf{X})}\right)}=n m \mathbb{E}\left[\|\nabla\| \cdot\left\|_{\ell_{p}^{n}(\mathbf{X})}\left(\frac{\mathrm{V}}{\| \mathrm{V} \ell_{\ell_{p}^{n}(\mathbf{X})}}\right)\right\|_{\ell_{2}^{n}\left(\ell_{2}^{m}\right)}\right] \tag{358}
\end{equation*}
$$

Also, using the formula from Lemma 155 for the density of $\|\mathrm{V}\|_{\ell_{p}^{n}(\mathbf{X})}$, for every $q>-n m$ we have

$$
\begin{equation*}
\mathbb{E}\left[\|\bigvee\|_{\ell_{p}^{n}(\mathbf{X})}^{q}\right]=\frac{\int_{0}^{\infty} s^{n m+q-1} e^{-s^{p}} \mathrm{~d} s}{\int_{0}^{\infty} r^{n m-1} e^{-r^{p}} \mathrm{~d} r}=\frac{\Gamma\left(\frac{n m+q}{p}\right)}{\Gamma\left(\frac{n m}{p}\right)} . \tag{359}
\end{equation*}
$$

Consequently,

$$
\begin{align*}
\mathbb{E}\left[\|\mathrm{V}\|_{\ell_{p}^{n}(\mathbf{X})}^{p-1}\|\nabla\| \cdot\left\|_{\ell_{p}^{n}(\mathbf{X})}\left(\frac{\mathrm{V}}{\|\mathrm{~V}\|_{\ell_{p}^{n}(\mathbf{X})}}\right)\right\|_{\ell_{2}^{n}\left(\ell_{2}^{m}\right)}\right] & =\mathbb{E}\left[\|\mathrm{V}\|_{\ell_{p}^{n}(\mathbf{X})}^{p-1}\right] \mathbb{E}\left[\|\nabla\| \cdot\left\|_{\ell_{p}^{n}(\mathbf{X})}\left(\frac{\mathrm{V}}{\|\mathrm{~V}\|_{\left.\ell_{p}^{n} \mathbf{X}\right)}}\right)\right\|_{\ell_{2}^{n}\left(\ell_{2}^{m}\right)}\right] \\
& =\frac{\Gamma\left(\frac{n m+p-1}{p}\right)}{n m \Gamma\left(\frac{n m}{p}\right)} \cdot \frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}(\mathbf{X})}\right.}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}(\mathbf{X})}\right)}, \tag{360}
\end{align*}
$$

where the first step of (360) uses the independence of $\|\mathrm{V}\|_{\ell_{p}^{n}(\mathbf{X})}$ and $\mathrm{V} /\|\mathrm{V}\|_{\ell_{p}^{n}(\mathbf{X})}$, by Lemma 155 , and the final step of (360) is a substitution of (358) and the case $q=p-1$ of (359). Hence,

$$
\begin{align*}
& \frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}(\mathbf{X})}\right)}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}(\mathbf{X})}\right)}=\frac{n m \Gamma\left(1+\frac{n m}{p}\right)}{\Gamma\left(1+\frac{n m-1}{p}\right)} \mathbb{E}\left[\|\mathrm{V}\|_{\ell_{p}^{n}(\mathbf{X})}^{p-1}\|\nabla\| \cdot \| \ell_{p}^{n}(\mathbf{X})\right. \\
&\left.\left(\frac{\mathrm{V}}{\|\mathrm{~V}\|_{\ell_{p}^{n}(\mathbf{X})}}\right) \|_{\ell_{2}^{n}\left(\ell_{2}^{m}\right.}\right]  \tag{361}\\
&=\frac{p \Gamma\left(1+\frac{n m}{p}\right)}{\Gamma\left(1+\frac{n m-1}{p}\right)} \mathbb{E}\left[\left(\sum_{i=1}^{n}\left\|\mathrm{~V}_{i}\right\|_{\mathbf{X}}^{2 p-2}\|\nabla\| \cdot\left\|\mathbf{X}\left(\frac{\mathrm{V}_{i}}{\left\|\mathrm{~V}_{i}\right\|_{\mathbf{X}}}\right)\right\|_{\ell_{2}^{m}}^{2}\right)^{\frac{1}{2}}\right],
\end{align*}
$$

where in the last step we used the identity (356).
The product structure of the density of V in (357) means that $\mathrm{V}_{1}, \ldots, \mathrm{~V}_{n}$ are (stochastically) independent. By Lemma 155 for each $i \in\{1, \ldots, n\}$ the random vector $\mathrm{V}_{i} /\left\|\mathrm{V}_{i}\right\|_{\mathbf{X}}$ is distributed on $\partial B_{\mathbf{X}}$ according to the cone measure $\kappa_{\mathbf{X}}$, and it is independent of the random variable

$$
\begin{equation*}
\mathrm{R}_{i} \stackrel{\text { def }}{=}\left\|\mathrm{V}_{i}\right\|_{\mathbf{X}}^{2 p-2}, \tag{362}
\end{equation*}
$$

whose density at $t \in(0, \infty)$ is equal (using Lemma 155 once more) to

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \operatorname{Prob}\left[\left\|\mathrm{~V}_{i}\right\|_{\mathrm{X}} \leqslant t^{\frac{1}{2 p-2}}\right]=\frac{\mathrm{d}}{\mathrm{~d} t} \int_{0}^{t^{\frac{1}{2 p-2}}} \frac{s^{m-1} e^{-s^{p}}}{\int_{0}^{\infty} r^{m-1} e^{-r^{p}} \mathrm{~d} r} \mathrm{~d} s=\frac{p}{2(p-1) \Gamma\left(\frac{m}{p}\right)} t^{\frac{m}{2 p-2}-1} e^{-t^{\frac{p}{2 p-2}}}
$$

Hence, the identity (361) which we established above coincides with the desired identity (354).
To prove the identity (355), let R be a random variable whose density at each $t \in(0, \infty)$ is given by (353), i.e., $\mathrm{R}_{1}, \ldots, \mathrm{R}_{n}$ are independent copies of R . Then, for every $\alpha>-m /(2 p-2)$ we have

$$
\begin{equation*}
\mathbb{E}\left[\mathrm{R}^{\alpha}\right]=\frac{p}{2(p-1) \Gamma\left(\frac{m}{p}\right)} \int_{0}^{\infty} t^{\frac{m}{2 p-2}+\alpha-1} e^{-\frac{p}{2 p-2}} \mathrm{~d} t=\frac{\Gamma\left(2 \alpha+\frac{m-2 \alpha}{p}\right)}{\Gamma\left(\frac{m}{p}\right)} . \tag{363}
\end{equation*}
$$

Using Lemma 155 (including the independence of $\mathrm{V}_{i} /\left\|\mathrm{V}_{i}\right\|_{\mathbf{X}}$ and $\left\|\mathrm{V}_{i}\right\|_{\mathbf{X}}$ ), we have

$$
\begin{align*}
\mathbb{E}\left[\sum_{i=1}^{n}\left\|\mathrm{~V}_{i}\right\|_{\mathbf{X}}^{2 p-2}\|\nabla\| \cdot\left\|_{\mathbf{X}}\left(\frac{\mathrm{V}_{i}}{\left\|\mathrm{~V}_{i}\right\|_{\mathbf{X}}}\right)\right\|_{\ell_{2}^{m}}^{2}\right] & =n \mathbb{E}[\mathrm{R}] \int_{\partial B_{\mathbf{X}}}\|\nabla\| \cdot\left\|_{\mathbf{X}}\right\|_{\ell_{2}^{m}}^{2} \mathrm{~d} \kappa_{\mathbf{X}} \\
& =\frac{n \Gamma\left(\frac{m+2 p-2}{p}\right)}{\Gamma\left(\frac{m}{p}\right)} \int_{\partial B_{\mathbf{X}}}\|\nabla\| \cdot\left\|_{\mathbf{X}}\right\|_{\ell_{2}^{m}}^{2} \mathrm{~d} \kappa_{\mathbf{X}}, \tag{364}
\end{align*}
$$

where we recall (362) and the last step of (364) is the case $\alpha=1$ of (363). At the same time,

$$
\begin{align*}
\mathbb{E}\left[\sum_{i=1}^{n}\left\|\mathrm{~V}_{i}\right\|_{\mathbf{X}}^{2 p-2}\|\nabla\| \cdot\left\|_{\mathbf{X}}\left(\frac{\mathrm{V}_{i}}{\left\|\mathrm{~V}_{i}\right\|_{\mathbf{X}}}\right)\right\|_{\ell_{2}^{m}}^{2}\right] & =\mathbb{E}\left[\|\mathrm{V}\|_{\ell_{p}^{n}(\mathbf{X})}^{2 p-2}\|\nabla\| \cdot\left\|_{\ell_{p}^{n}(\mathbf{X})}\left(\frac{\mathrm{V}}{\|\mathrm{~V}\|_{\ell_{( }^{n}(\mathbf{X})}}\right)\right\|_{\ell_{2}^{n}\left(\ell_{2}^{m}\right)}^{2}\right] \\
& =\mathbb{E}\left[\|\mathrm{V}\|_{\ell_{p}^{n}(\mathbf{X})}^{2 p-2}\right] \mathbb{E}\left[\|\nabla\| \cdot\left\|_{\ell_{p}^{n}(\mathbf{X})}\left(\frac{\mathrm{V}}{\|\mathrm{~V}\|_{\ell_{p}^{n}(\mathbf{X})}}\right)\right\|_{\ell_{2}^{n}\left(\ell_{2}^{m}\right)}^{2}\right]  \tag{365}\\
& =\frac{\Gamma\left(\frac{n m+2 p-2}{p}\right)}{\Gamma\left(\frac{n m}{p}\right)} \int_{\partial B_{\ell_{p}^{n}(\mathbf{X})}}\|\nabla\| \cdot\left\|_{\ell_{p}^{n}(\mathbf{X})}\right\|_{\ell_{2}^{n}\left(\ell_{2}^{m}\right)}^{2} \mathrm{~d} \kappa_{\ell_{p}^{n}(\mathbf{X})},
\end{align*}
$$

where the first step of (365) uses the identity (356), the second step of (365) uses the independence of $\|\mathrm{V}\|_{\ell_{p}^{n}(\mathbf{X})}$ and $\mathrm{V} /\|\mathrm{V}\|_{\ell_{p}^{n}(\mathbf{X})}$ per Lemma 155 and the final step of uses the case $q=2 p-2$ of (359) and Lemma 155. The desired identity (355) now follows by substituting (365) into (364).

The following lemma will have a central role in the proof of Theorem 24 and Theorem 47 ,
Lemma 157. Suppose that $n, m \in \mathbb{N}$ and $\beta>0$ satisfy $\beta \leqslant \frac{m-1}{2}$. Then, for every $1 \leqslant p \leqslant m$ we have

$$
\mathrm{iq}\left(B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)=\sqrt{n m}=\sqrt{\operatorname{dim}\left(\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)\right)},
$$

where we recall that the normed space $\Omega_{\beta}^{m}=\left(\mathbb{R}^{m},\|\cdot\|_{\Omega_{\beta}^{m}}\right.$ was defined in (342).
Prior to proving Lemma 157 , we will show how it implies Theorem 47 , and then deduce Theorem 24
Proof of Theorem 47 assuming Lemma 157. By the assumption 47) of Theorem 47, write $n=k m$ for some $k, m \in \mathbb{N}$ with $\max \{2, p\} \leqslant m \leqslant e^{p}$. Then $(m-1) / 2>0$ and $m \geqslant p$, so we may apply Lemma 157 with $n$ replaced by $k$ and $\beta=(m-1) / 2$. Denoting $\mathbf{Y}=\ell_{p}^{k}\left(\Omega_{\beta}^{m}\right)$, the conclusion of Lemma 157 is that $\mathrm{iq}\left(B_{\mathbf{Y}}\right)=\sqrt{n}$.
$\mathbf{Y}$ is canonically positioned (it is a space from Example 41). To prove Theorem 47, it remains to check that $\|\cdot\|_{\mathbf{Y}}=\|\cdot\|_{\ell_{p}^{n}}$, where, since $n=k m$, we identify $\mathbb{R}^{n}$ with $\mathrm{M}_{k \times n}(\mathbb{R})$, namely we identify $\ell_{p}^{n}$ with $\ell_{p}^{k}\left(\ell_{p}^{m}\right)$.

In fact, for any $\beta>0$ (not only our choice $\beta=(m-1) / 2$ above) we will check that

$$
\begin{equation*}
\forall x \in \mathbb{R}^{m}, \quad\left(1-e^{-\frac{\beta}{m}}\right)\|x\|_{\Omega_{\beta}^{m}} \leqslant\|x\|_{\ell_{\infty}^{m}} \leqslant\|x\|_{\Omega_{\beta}^{m}} . \tag{366}
\end{equation*}
$$

It follows from (366) that $\|\cdot\|_{\Omega_{\beta}^{m}}=\|\cdot\|_{\ell_{\infty}^{m}}$ when $\beta=m$. But, $\|\cdot\|_{\ell_{p}^{m}}=\|\cdot\|_{\ell_{\infty}^{m}}$ by the assumption $e^{p} \geqslant m$. So,

$$
\beta=n \Longrightarrow\|\cdot\|_{\mathbf{Y}}=\|\cdot\|_{\ell_{p}^{k}\left(\Omega_{\beta}^{m}\right)}=\|\cdot\|_{\ell_{p}^{k}\left(\ell_{\infty}^{m}\right)}=\|\cdot\|_{\ell_{p}^{k}\left(\ell_{p}^{m}\right)}=\|\cdot\|_{\rho}^{n} .
$$

Fix $x \in \mathbb{R}^{m}$. To verify the second inequality in (366), the definition (342) gives $\sum_{i=1}^{m} \psi_{\beta}\left(\left|x_{i}\right| / s\right)=\infty$ when $0<s \leqslant\|x\|_{\ell_{\infty}^{m}}$, so $\|x\|_{\Omega_{\beta}^{m}} \geqslant\|x\|_{\ell_{\infty}^{m}}$ by (340). For the first inequality in (366), by direct differentiation it is elementary to verify that the function $u \mapsto \log (1 /(1-u)) / u$ is increasing on the interval $[0,1)$. Thus,

$$
0 \leqslant t \leqslant \alpha<1 \Longrightarrow \psi_{\beta}(t)=\frac{1}{\beta} \log \left(\frac{1}{1-t}\right) \leqslant \frac{\log \left(\frac{1}{1-\alpha}\right)}{\alpha \beta} t .
$$

Hence, for every fixed $0<\alpha<1$,

$$
\begin{equation*}
s \geqslant \frac{1}{\alpha}\|x\|_{\ell_{\infty}^{m}} \Longrightarrow \sum_{i=1}^{m} \psi_{\beta}\left(\frac{\left|x_{i}\right|}{s}\right) \leqslant \sum_{i=1}^{m} \frac{\log \left(\frac{1}{1-\alpha}\right)}{\alpha \beta s}\left|x_{i}\right| \leqslant \frac{m \log \left(\frac{1}{1-\alpha}\right)}{\alpha \beta s}\|x\|_{\ell_{\infty}^{m}} . \tag{367}
\end{equation*}
$$

Provided $\alpha \geqslant 1-e^{-\beta / m}$, the choice $s=m \log (1 /(1-\alpha))\|x\|_{\ell_{\infty}^{m}} /(\alpha \beta)$ satisfies the requirement $s \geqslant\|x\|_{\ell_{\infty}^{m}} / \alpha$, so we get from (340) and (367) that

$$
\begin{equation*}
\|x\|_{\Omega_{\beta}^{m}} \leqslant \frac{m \log \left(\frac{1}{1-\alpha}\right)}{\alpha \beta}\|x\|_{\ell_{\infty}^{m}} . \tag{368}
\end{equation*}
$$

The optimal choice of $\alpha$ in (368) is $\alpha=1-e^{-\beta / m}$, giving the first inequality in (366).
Having proved Theorem 47 (assuming Lemma 157, which we will soon prove), we have also already established Theorem 24 provided $n \in \mathbb{N}$ and $p \geqslant 1$ satisfy the divisor condition (74). Indeed, the space $\mathbf{Y}$ that Theorem 47 provides is canonically positioned and hence by the discussion in Section 1.6 .2 it is also in its minimum surface area position, so by [GP99, Proposition 3.1] we have

$$
\frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)}=\frac{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right) \sqrt{n}}=\left(\frac{\mathrm{iq}\left(B_{\mathbf{Y}}\right)}{\sqrt{n}}\right) \frac{1}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)^{\frac{1}{n}}}=\frac{1}{\operatorname{vol}_{n}\left(B_{\ell_{p}^{n}}\right)^{\frac{1}{n}}} \stackrel{\sqrt[315]{=}}{\stackrel{\frac{1}{p}}{p},}
$$

where the penultimate step uses the fact that $\mathrm{iq}\left(B_{\mathbf{Y}}\right)=\sqrt{n}$ by Theorem 47 , and also that by Theorem 47 we have $\|\cdot\|_{\mathbf{Y}}=\|\cdot\|_{\ell_{p}^{n}}$, which implies that the $n^{\prime}$ th root of the volume of the unit ball of $\mathbf{Y}$ is proportional to the $n$ 'th root of the volume of the unit ball of $\ell_{p}^{n}$.

The deduction of Theorem 24 for the remaining values of $p \geqslant 1$ and $n \in \mathbb{N}$ uses the following identity, which we will also use in the proof of Proposition 163 below.

Lemma 158. Fix $n, m \in \mathbb{N}$. Suppose that $K \subseteq \mathbb{R}^{n}$ and $L \subseteq \mathbb{R}^{m}$ are convex bodies. Then,

$$
\frac{\operatorname{MaxProj}(K \times L)}{\operatorname{vol}_{n+m}(K \times L)}=\left(\frac{\operatorname{MaxProj}(K)^{2}}{\operatorname{vol}_{n}(K)^{2}}+\frac{\operatorname{MaxProj}(L)^{2}}{\operatorname{vol}_{m}(L)^{2}}\right)^{\frac{1}{2}} .
$$

Proof. Fix $z \in S^{n+m-1}$. By the Cauchy projection formula Gar06] that we recalled in 30, we have

$$
\operatorname{vol}_{n+m-1}\left(\operatorname{Proj}_{z^{\perp}}(K \times L)\right)=\frac{1}{2} \int_{\partial(K \times L)}\left|\left\langle z, N_{K \times L}(w)\right\rangle\right| \mathrm{d} w,
$$

where $N_{K \times L}(w)$ is the (almost-everywhere defined) unit outer normal to $\partial(K \times L)$ at $w \in \partial(K \times L)$. Now,

$$
\partial(K \times L)=(\partial K \times L) \cup(K \times \partial L) \quad \text { and } \quad \operatorname{vol}_{n+m-1}((\partial K \times L) \cap(K \times \partial L))=0
$$

Consequently,

$$
\operatorname{vol}_{n+m-1}\left(\operatorname{Proj}_{z^{\perp}}(K \times L)\right)=\frac{1}{2} \int_{\partial K \times L}\left|\left\langle z, N_{K \times L}(w)\right\rangle\right| \mathrm{d} w+\frac{1}{2} \int_{K \times \partial L}\left|\left\langle z, N_{K \times L}(w)\right\rangle\right| \mathrm{d} w
$$

If we write each $w \in \mathbb{R}^{n}$ as $w=\left(w_{1}, w_{2}\right)$ where $w_{1} \in \mathbb{R}^{n}$ and $w_{2} \in \mathbb{R}^{m}$, then for almost every (with respect to the $(n+m-1)$-dimensional Hausdorff measure) $w \in \partial K \times L$ we have $N_{K \times L}(w)=\left(N_{K}\left(w_{1}\right), 0\right)$. Also, for almost every $w \in K \times \partial L$ we have $N_{K \times L}(w)=\left(0, N_{L}\left(w_{2}\right)\right)$. We therefore have

$$
\begin{aligned}
\operatorname{vol}_{n+m-1}\left(\operatorname{Proj}_{z^{\perp}}(K \times L)\right) & =\frac{\operatorname{vol}_{m}(L)}{2} \int_{\partial K}\left|\left\langle z_{1}, N_{K}(x)\right\rangle\right| \mathrm{d} x+\frac{\operatorname{vol}_{n}(K)}{2} \int_{\partial L}\left|\left\langle z_{2}, N_{L}(y)\right\rangle\right| \mathrm{d} y \\
& =\operatorname{vol}_{m}(L) \operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z_{1}^{\perp}} K\right)\left\|z_{1}\right\|_{\ell_{2}^{n}}+\operatorname{vol}_{n}(K) \operatorname{vol}_{m-1}\left(\operatorname{Proj}_{z_{2}^{\perp}} L\right)\left\|z_{2}\right\|_{\ell_{2}^{m}}
\end{aligned}
$$

where the last step is two applications of the Cauchy projection formula (in $\mathbb{R}^{n}$ and $\mathbb{R}^{m}$ ). Hence,

$$
\begin{aligned}
\frac{\operatorname{vol}_{n+m-1}\left(\operatorname{Proj}_{z^{\perp}}(K \times L)\right)}{\operatorname{vol}_{n+m}(K \times L)} & =\frac{\operatorname{vol}_{n+m-1}\left(\operatorname{Proj}_{z^{\perp}}(K \times L)\right)}{\operatorname{vol}_{n}(K) \operatorname{vol}_{m}(L)} \\
& =\frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{z_{1}^{\perp}} K\right)}{\operatorname{vol}_{n}(K)}\left\|z_{1}\right\|_{\ell_{2}^{n}}+\frac{\operatorname{vol}_{m-1}\left(\operatorname{Proj}_{z_{2}^{\perp}} L\right)}{\operatorname{vol}_{m}(L)}\left\|z_{2}\right\|_{\ell_{2}^{m}} .
\end{aligned}
$$

Consequently,

$$
\begin{aligned}
\frac{\operatorname{MaxProj}(K \times L)}{\operatorname{vol}_{n+m}(K \times L)} & =\max _{z \in S^{n+m-1}} \frac{\operatorname{vol}_{n+m-1}\left(\operatorname{Proj}_{z^{\perp}}(K \times L)\right)}{\operatorname{vol}_{n+m}(K \times L)} \\
& =\max _{(u, v) \in S^{1}} \max _{x \in S^{n-1}} \max _{y \in S^{m-1}} \frac{\operatorname{vol}_{n+m-1}\left(\operatorname{Proj}_{(u x+v y)^{\perp}}(K \times L)\right)}{\operatorname{vol}_{n+m}(K \times L)} \\
& =\max _{(u, v) \in S^{1}} \max _{x \in S^{n-1}} \max _{y \in S^{m-1}}\left(\frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}} K\right)}{\operatorname{vol}_{n}(K)}|u|+\frac{\operatorname{vol}_{m-1}\left(\operatorname{Proj}_{y^{\perp}} L\right)}{\operatorname{vol}_{m}(L)}|v|\right) \\
& =\max _{(u, v) \in S^{1}}\left(\frac{\operatorname{MaxProj}(K)}{\operatorname{vol}_{n}(K)}|u|+\frac{\operatorname{MaxProj}(L)}{\operatorname{vol}_{m}(L)}|v|\right) \\
& =\left(\frac{\operatorname{MaxProj}(K)^{2}}{\operatorname{vol}_{n}(K)^{2}}+\frac{\operatorname{MaxProj}(L)^{2}}{\operatorname{vol}_{m}(L)^{2}}\right)^{\frac{1}{2}} .
\end{aligned}
$$

We can now prove Theorem 24 in its full generality using the fact that we proved Theorem 47,
Proof of Theorem 24. Let $m$ be any integer that satisfies $\max \{2, p\} \leqslant m \leqslant e^{p}$ (if $1 \leqslant p \leqslant 2$, then take $m=2$, and if $p \geqslant 2$, then such an $m$ exists because $e^{p}-p \geqslant e^{2}-2>5$ ). Write $n=k m+r$ for some $k \in \mathbb{N} \cup\{0\}$ and $r \in\{0, \ldots, m-1\}$. If $r=0$, then $m$ divides $n$ and we can conclude by applying Theorem47as we did above (recall the paragraph immediately before Lemma 158). So, assume from now that $r \geqslant 1$.

By Theorem 47 there is a canonically positioned normed space $\mathbf{Y}=\left(\mathbb{R}^{k m},\|\cdot\|_{\mathbf{Y}}\right)$ such that $\mathrm{iq}\left(B_{\mathbf{Y}}\right)=\sqrt{\mathrm{km}}$ and $\|\cdot\|_{\mathbf{Y}}=\|\cdot\|_{\ell_{p}^{k m}}$. Define $\mathbf{Y}_{p}^{n}=\mathbf{Y} \oplus_{\infty} \Omega_{\beta}^{r}$, where $\beta=r$ and iq $\left(\Omega_{\beta}^{r}\right)=\sqrt{r}$; such $\beta$ exists trivially if $r=1$, and if $r \geqslant 2$, then its existence follows from an application of Lemma 157 (with the choices $n=1$ and $p=m=r$ ).

Since $\beta=r$, by 366) we have $\|\cdot\|_{\Omega_{\beta}^{r}}=\|\cdot\|_{\ell_{\infty}^{r}}$. Also, $\|\cdot\|_{\ell_{\infty}^{r}}=\|\cdot\|_{\ell_{p}^{r}}$ since $e^{p} \geqslant m>r$. Consequently,

$$
\forall(x, y) \in \mathbb{R}^{k m} \times \mathbb{R}^{r}, \quad \max \left\{\|x\|_{\mathbf{Y}},\|y\|_{\Omega_{\beta}^{r}}\right\}=\max \left\{\|x\|_{\ell_{p}^{k m}},\|y\|_{\ell_{p}^{r}}\right\}=\left(\|x\|_{\ell_{p}^{k m}}^{p}+\|y\|_{\ell_{p}^{r}}^{p}\right)^{\frac{1}{p}}
$$

Recalling the definition of $\mathbf{Y}_{p}^{n}$, this means that $\|\cdot\|_{\mathbf{Y}_{p}^{n}}=\|\cdot\|_{\ell_{p}^{n}}$.
Since both $\mathbf{Y}$ and $\Omega_{\beta}^{r}$ are canonically positioned and hence in their minimum surface area positions,

$$
\frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}}\right)}{\operatorname{vol}_{k m}\left(B_{\mathbf{Y}}\right)}=\left(\frac{\mathrm{iq}\left(B_{\mathbf{Y}}\right)}{\sqrt{k m}}\right) \frac{1}{\operatorname{vol}_{k m}\left(B_{\mathbf{Y}}\right)^{\frac{1}{k m}}}=\frac{1}{\operatorname{vol}_{k m}\left(B_{\ell_{p}^{k m}}\right)^{\frac{1}{k m}}}=(k m)^{\frac{1}{p}},
$$

and

$$
\frac{\operatorname{MaxProj}\left(B_{\Omega_{\beta}^{r}}\right)}{\operatorname{vol}_{r}\left(B_{\Omega_{\beta}^{r}}\right)}=\left(\frac{\operatorname{iq}\left(\Omega_{\beta}^{r}\right)}{\sqrt{r}}\right) \frac{1}{\operatorname{vol}\left(\Omega_{\beta}^{r}\right)^{\frac{1}{r}}}=\frac{1}{\operatorname{vol}\left(\ell_{\infty}^{r}\right)^{\frac{1}{r}}}=1=r^{\frac{1}{p}} .
$$

Consequently, since $B_{\mathbf{Y}_{p}^{n}}=B_{\mathbf{Y}} \times B_{\Omega_{\beta}^{r}}$, by Lemma 158 we conclude that

$$
\frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}_{p}^{n}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}_{p}^{n}}\right)}=\left(\frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}}\right)^{2}}{\operatorname{vol}_{k m}\left(B_{\mathbf{Y}}\right)^{2}}+\frac{\operatorname{MaxProj}\left(B_{\Omega_{\beta}^{r}}\right)^{2}}{\operatorname{vol}_{r}\left(B_{\Omega_{\beta}^{r}}\right)^{2}}\right)^{\frac{1}{2}}=\left((k m)^{\frac{2}{p}}+r^{\frac{2}{p}}\right)^{\frac{1}{2}}=(k m+r)^{\frac{1}{p}}=n^{\frac{1}{p}}
$$

The following lemma will be used in the proof of Lemma 157
Lemma 159. Suppose that $m \in \mathbb{N}, r \in \mathbb{N} \cup\{0\}$ and $\beta>0$ satisfy $\beta \leqslant \frac{m+r-2}{2}$. Then

$$
\begin{equation*}
\int_{\partial B_{\ell_{1}^{m}}}\left(e^{\beta\left|\tau_{1}\right|}-\sum_{k=r-1}^{\infty} \frac{\beta^{k}\left|\tau_{1}\right|^{k}}{k!}\right) \mathrm{d} \kappa_{\ell_{1}^{m}}(\tau)=\int_{\partial B_{\ell_{1}^{m}}}\left(\sum_{k=r}^{\infty} \frac{\beta^{k}\left|\tau_{1}\right|^{k}}{k!}\right) \mathrm{d} \kappa_{\ell_{1}^{m}}(\tau)=\frac{\beta^{r}(m-1)!}{(m+r-1)!} . \tag{369}
\end{equation*}
$$

Proof. Let $\mathrm{H}_{1}, \ldots, \mathrm{H}_{m}$ be independent random variables whose density at each $s \in \mathbb{R}$ is equal to $e^{-|s|} / 2$. Then, $\left|\mathrm{H}_{1}\right|, \ldots,\left|H_{m}\right|$ are exponential random variables of rate 1 , and therefore if we denote

$$
\Gamma \stackrel{\text { def }}{=} \sum_{i=1}^{m}\left|\mathrm{H}_{i}\right|,
$$

then $\Gamma$ has $\Gamma(m, 1)$ distribution, i.e., its density at each $s \geqslant 0$ is equal to $s^{m-1} e^{-s} /(m-1)!$; the proof of this standard probabilistic fact can be found in e.g. Dur19]. By [SZ90, RR91a] (or Lemma 155], the random vector $\left(\mathrm{H}_{1}, \ldots, \mathrm{H}_{m}\right) / \Gamma$ is distributed according to $\kappa_{\ell_{1}^{m}}$ and is independent of $\Gamma$. Thus, for every $k \in \mathbb{N}$,

$$
\int_{\partial B_{e_{1}^{m}}}\left|\tau_{1}\right|^{k} \mathrm{~d} \kappa_{\ell_{1}^{m}}(\tau)=\mathbb{E}\left[\frac{\left|\mathrm{H}_{1}\right|^{k}}{\Gamma^{k}}\right]=\frac{\mathbb{E}\left[\left|\mathrm{H}_{1}\right|^{k}\right]}{\mathbb{E}\left[\Gamma^{k}\right]}=\frac{\int_{0}^{\infty} s^{k} e^{-s} \mathrm{~d} s}{\frac{1}{(m-1)!} \int_{0}^{\infty} s^{k+m-1} e^{-s} \mathrm{~d} s}=\frac{k!(m-1)!}{(k+m-1)!} .
$$

Consequently,

$$
\begin{equation*}
\int_{\partial B_{\ell_{1}^{m}}}\left(\sum_{k=r}^{\infty} \frac{\beta^{k}\left|\tau_{1}\right|^{k}}{k!}\right) \mathrm{d} \kappa_{\ell_{1}^{m}}(\tau)=\frac{(m-1)!}{\beta^{m-1}} \sum_{k=r}^{\infty} \frac{\beta^{k+m-1}}{(k+m-1)!}=\frac{\beta^{r}(m-1)!}{(m+r-2)!} \int_{0}^{1} e^{\beta t}(1-t)^{m+r-2} \mathrm{~d} t, \tag{370}
\end{equation*}
$$

where the last step is the integral form of the remainder of the Taylor series of the exponential function.
It is mechanical to check that (369) holds for $m \in\{1,2\}$, so assume for the rest of the proof of Lemma 159 that $m \geqslant 3$. We then see from (370) that our goal (369) is equivalent to showing that

$$
\begin{equation*}
\int_{0}^{1} e^{\beta t}(1-t)^{m+r-2} \mathrm{~d} t=\frac{1}{m+r} . \tag{371}
\end{equation*}
$$

For the upper bound in (371), estimate the integrand using $(1-t)^{m+r-2} \leqslant e^{-(m+r-2) t}$ to get

$$
\int_{0}^{1} e^{\beta t}(1-t)^{m+r-2} \mathrm{~d} t \leqslant \int_{0}^{1} e^{-(m+r-2-\beta) t} \mathrm{~d} t=\frac{1-e^{-(m+r-2-\beta)}}{m+r-2-\beta}=\frac{1}{m+r},
$$

where we used $\beta<\frac{m+r-2}{2}$. For the lower bound in (371), since $(1-t)^{m+r-2} \gtrsim 1$ when $0 \leqslant t \leqslant \frac{1}{m+r-2}$,

$$
\int_{0}^{1} e^{\beta t}(1-t)^{m+r-2} \mathrm{~d} t \geqslant \int_{0}^{\frac{1}{m+r-2}} e^{\beta t}(1-t)^{m+r-2} \mathrm{~d} t \gtrsim \int_{0}^{\frac{1}{m+r-2}} e^{\beta t} \mathrm{~d} t=\frac{e^{\frac{\beta}{m+r-2}}-1}{\beta}=\frac{1}{m+r},
$$

where in the last step we used the assumption $\beta<\frac{m+r-2}{2}$ once more.
Proof of Lemma 157. By combining the case $g \equiv 1$ of (341) with (344), we see that

$$
\begin{equation*}
\operatorname{vol}_{m}\left(B_{\Omega_{\beta}^{m}}\right)=\frac{\beta^{m-1} 2^{m}}{e^{\beta} m!} m \int_{\partial B_{\ell_{1}^{m}}}\left(e^{\beta\left|\tau_{1}\right|}-1\right) \mathrm{d} \kappa_{\ell_{1}^{m}}(\tau) \stackrel{\sqrt{369}}{=} \frac{(2 \beta)^{m}}{e^{\beta} m!} . \tag{372}
\end{equation*}
$$

Since we are assuming in Lemma 157 that $\beta \lesssim m$, in combination with (315) we get from (372) that

$$
\begin{equation*}
\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)^{\frac{1}{n m}}=\frac{\beta}{n^{\frac{1}{p}} m} \tag{373}
\end{equation*}
$$

At the same time, by applying Cauchy-Schwarz to the identity (354) of Lemma 156 we have

$$
\begin{align*}
\frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)} & \leqslant \frac{p \Gamma\left(1+\frac{n m}{p}\right)}{\Gamma\left(1+\frac{n m-1}{p}\right)}\left(n\left(\mathbb{E}\left[\mathrm{R}_{1}\right]\right) \int_{\partial B_{\Omega_{\beta}^{m}}}\|\nabla\| \cdot\left\|_{\Omega_{\beta}^{m}}(\theta)\right\|_{\ell_{2}^{m}}^{2} \mathrm{~d} \kappa_{\Omega_{\beta}^{m}}(\theta)\right)^{\frac{1}{2}}  \tag{374}\\
& =n^{\frac{1}{p}+\frac{1}{2}} m\left(\int_{\partial{\Omega_{\Omega}^{m}}}\|\nabla\| \cdot\left\|_{\Omega_{\beta}^{m}}(\theta)\right\|_{\ell_{2}^{m}}^{2} \mathrm{~d} \kappa_{\Omega_{\beta}^{m}}(\theta)\right)^{\frac{1}{2}}
\end{align*}
$$

where the random variable $R_{1}$ is as in Lemma 156, i.e., its density is in (353), and the last step is an application the evaluation (363) of its moments and Stirling's formula, using the assumption $1 \leqslant p \leqslant m$.

Recalling (342), even though $\|\cdot\|_{\Omega_{\beta}^{m}}$ is defined implicitly by (340), we can compute $\nabla\|\cdot\|_{\Omega_{\beta}^{m}}(\theta)$ for almost every $\theta \in \partial B_{\Omega_{\beta}^{m}}$ as the unique vector $v \in \mathbb{R}^{m}$ that is normal to $\partial B_{\Omega_{\beta}^{m}}$ and satisfies $\langle v, \theta\rangle=1$. Indeed, since $\partial \Omega_{\beta}^{m}$ is parameterized as the zero set of the function $\Psi_{\beta}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ that is given by

$$
\forall x \in \mathbb{R}^{n}, \quad \Psi_{\beta}(x) \stackrel{\text { def }}{=} 1-\sum_{i=1}^{m} \psi_{\beta}\left(\left|x_{i}\right|\right),
$$

the following vector is normal to $\partial B_{\Omega_{\beta}^{m}}$ for almost every $\theta \in \partial B_{\Omega_{\beta}^{m}}$.

$$
v_{\beta}(\theta) \stackrel{\text { def }}{=} \nabla \Psi_{\beta}(\theta)=-\left(\psi_{\beta}^{\prime}\left(\left|\theta_{1}\right|\right) \operatorname{sign}\left(\theta_{1}\right), \ldots, \psi_{\beta}^{\prime}\left(\left|\theta_{m}\right|\right) \operatorname{sign}\left(\theta_{m}\right)\right) .
$$

So, $\nabla\|\cdot\|_{\Omega_{\beta}^{m}}(\theta)=\lambda_{\beta}(\theta) v_{\beta}(\theta)$ for almost every $\theta \in \partial B_{\Omega_{\beta}^{m}}$, where $\lambda_{\beta}(\theta) \in \mathbb{R}$ is such that $\left\langle\lambda_{\beta}(\theta) v_{\beta}(\theta), \theta\right\rangle=1$, i.e., $\lambda_{\beta}(\theta)=-1 /\left\langle v_{\beta}(\theta), \theta\right\rangle$. This shows that for almost every $\theta \in \partial B_{\Omega_{\beta}^{m}}$,

$$
\begin{align*}
\nabla\|\cdot\|_{\Omega_{\beta}^{m}}(\theta) & =\frac{1}{\sum_{i=1}^{m}\left|\theta_{i}\right| \psi_{\beta}^{\prime}\left(\left|\theta_{i}\right|\right)}\left(\psi_{\beta}^{\prime}\left(\left|\theta_{1}\right|\right) \operatorname{sign}\left(\theta_{1}\right), \ldots, \psi_{\beta}^{\prime}\left(\left|\theta_{m}\right|\right) \operatorname{sign}\left(\theta_{m}\right)\right) \\
& =\frac{1}{\sum_{i=1}^{m} \frac{\left|\theta_{i}\right|}{1-\left|\theta_{i}\right|}}\left(\frac{\operatorname{sign}\left(\theta_{1}\right)}{1-\left|\theta_{1}\right|}, \ldots, \frac{\operatorname{sign}\left(\theta_{m}\right)}{1-\left|\theta_{m}\right|}\right), \tag{375}
\end{align*}
$$

where the first equality in (375) holds for any $\psi_{\beta}$ that satisfies the conditions of Lemma 153 , and for the second equality in (375) recall the definition (342) of the specific $\psi_{\beta}$ that we are using here. Therefore,

$$
\begin{align*}
& \int_{\partial B_{\Omega_{\beta}^{m}}}\|\nabla\| \cdot\left\|_{\Omega_{\beta}^{m}}(\theta)\right\|_{\ell_{2}^{m}}^{2} \mathrm{~d} \kappa_{\Omega_{\beta}^{m}}(\theta)=\frac{\left.\int_{\partial B_{\ell_{1}^{m}}} \frac{\sum_{i=1}^{m}}{\sum_{i=1}^{m}\left(e^{\beta\left|\tau_{i}\right|}|1|\right.} \mathrm{d}\right)}{} \mathrm{d} \kappa_{\ell_{1}^{m}}(\tau) \\
& m \int_{\partial B_{\ell_{1}^{m}}}\left(e^{\beta\left|\tau_{1}\right|}-1\right) \mathrm{d} \kappa_{\ell_{1}^{m}}^{m}(\tau)  \tag{376}\\
& \leqslant \frac{\int_{\partial B_{\ell_{1}^{m}}} \frac{\sum_{i=1}^{m} e^{2 \beta\left|\tau_{i}\right|}}{\beta \sum_{i=1}^{m}\left|\tau_{i}\right|} \mathrm{d} \kappa_{\ell_{1}^{m}}(\tau)}{m \int_{\partial B_{\ell_{1}^{m}}}\left(e^{\beta\left|\tau_{1}\right|}-1\right) \mathrm{d} \kappa_{\ell_{1}^{m}}^{m}(\tau)}=\frac{\int_{\partial B_{\ell_{1}^{m}}} e^{2 \beta\left|\tau_{1}\right|} \mathrm{d} \kappa_{\ell_{1}^{m}}(\tau)}{\beta \int_{\partial B_{\ell_{1}^{m}}}\left(e^{\beta\left|\tau_{1}\right|}-1\right) \mathrm{d} \kappa_{\ell_{1}^{m}(\tau)}}=\frac{m}{\beta^{2}},
\end{align*}
$$

where the first step of (376) is a substitution of (375) into (343) while using (344) and that $\psi_{\beta}^{-1}(t)=1-e^{-\beta t}$ for every $t \geqslant 0$, the second step of (376) uses the inequality $e^{t} \geqslant t+1$ which holds for any $t \in \mathbb{R}$, and the final step of (376) is an application of Lemma 159. Now, a combination of (374) and (376) gives

$$
\begin{equation*}
\frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)} \lesssim \frac{n^{\frac{1}{p}+\frac{1}{2}} m^{\frac{3}{2}}}{\beta} \tag{377}
\end{equation*}
$$

By combining (373) and (377) we conclude that

$$
\operatorname{iq}\left(B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)=\frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)} \operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}\left(\Omega_{\beta}^{n}\right)}\right)^{\frac{1}{n m}} \lesssim \sqrt{n m} .
$$

The reverse inequality $\operatorname{iq}\left(B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right) \gtrsim \sqrt{n m}$ follows from the isoperimetric theorem (12), so the proof of Lemma 157 is complete. Note that this also shows that all of the inequalities that we derived in the above proof of Lemma 157 are in fact asymptotic equivalences. This holds in particular for 377, i.e.,

$$
\frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}\left(\Omega_{\beta}^{m}\right)}\right)}=\frac{n^{\frac{1}{p}+\frac{1}{2}} m^{\frac{3}{2}}}{\beta} .
$$

The following asymptotic evaluation of the surface area of the sphere of $\ell_{p}^{n}\left(\ell_{q}^{m}\right)$ in the entire range of possible values of $p, q \geqslant 1$ and $m, n \in \mathbb{N}$ is an application of Lemma 156; by (316) it is equivalent to (83).

Theorem 160. For every $n, m \in \mathbb{N}$ and $p, q \in[1, \infty]$ we have

$$
\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}\left(e_{q}^{m}\right)}\right)=\frac{2^{n m} \Gamma\left(1+\frac{1}{q}\right)^{n m} \Gamma\left(1+\frac{m}{p}\right)^{n}}{\Gamma\left(1+\frac{m}{q}\right)^{n} \Gamma\left(1+\frac{n m}{p}\right)} \cdot \begin{cases}n^{1+\frac{1}{p}} m^{1+\frac{1}{q}} & m \leqslant \min \left\{\frac{p}{n}, q\right\},  \tag{378}\\ \sqrt{q} n^{1+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} & q \leqslant m \leqslant \frac{p}{n}, \\ \sqrt{p} n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} & \frac{p}{n} \leqslant m \leqslant \min \{p, q\}, \\ \sqrt{p q} n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{q}} & \max \left\{\frac{p}{n}, q\right\} \leqslant m \leqslant p, \\ n^{\frac{1}{2}+\frac{1}{p}} m^{1+\frac{1}{q}} & p \leqslant m \leqslant q, \\ \sqrt{q} n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} & m \geqslant \max \{p, q\} .\end{cases}
$$

Proof. By continuity we may assume that $p, q \in(1, \infty)$. Suppose that G is a symmetric real-valued random variable whose density at each $s \in \mathbb{R}$ is equal to

$$
\begin{equation*}
\frac{1}{2 \Gamma\left(1+\frac{1}{q}\right)} e^{-|s|^{q}} . \tag{379}
\end{equation*}
$$

Let $\mathrm{G}_{1}, \ldots, \mathrm{G}_{m}$ be independent copies of G . Set $\mathrm{U} \stackrel{\text { def }}{=}\left(\mathrm{G}_{1}, \ldots, \mathrm{G}_{m}\right) \in \mathbb{R}^{m}$. By the probabilistic representation of the cone measure on $\partial B_{\ell_{q}^{m}}$ in SZ90, RR91a] (or Lemma 155), the random vector $\mathrm{U} /\|\mathrm{U}\|_{\ell_{q}^{m}}$ is distributed according to the cone measure on $\partial B_{\ell_{q}^{m}}$, and moreover it is independent of $\|\mathrm{U}\|_{\ell_{q}^{m}}$.

Consider the following random variable.

$$
\begin{equation*}
\mathrm{N} \stackrel{\text { def }}{=}\|\nabla\| \cdot\left\|_{\ell_{q}^{m}}\left(\frac{\mathrm{U}}{\|\mathrm{U}\|_{\ell_{q}^{m}}}\right)\right\|_{\ell_{2}^{m}}^{2}=\frac{1}{\|\mathrm{U}\|_{\ell_{q}^{m}}^{2-2}} \sum_{j=1}^{m}\left|\mathrm{G}_{j}\right|^{2 q-2}=\frac{\|\mathrm{U}\|_{\ell_{2 q-2}^{m}}^{2 q-2}}{\|\mathrm{U}\|_{\ell_{q}^{m}}^{2 q-2}} . \tag{380}
\end{equation*}
$$

If we let $N_{1}, \ldots, N_{n}, \mathrm{R}_{1}, \ldots, \mathrm{R}_{n}$ be independent random variables such that $\mathrm{N}_{1}, \ldots, \mathrm{~N}_{n}$ have the same distribution as $N$, and $R_{1}, \ldots, R_{n}$ are as in Lemma 156, then by Lemma 156] we have

$$
\begin{equation*}
\frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}\left(\ell_{q}^{m}\right)}\right)}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}\left(\ell_{q}^{m}\right)}\right)}=\frac{p \Gamma\left(1+\frac{n m}{p}\right)}{\Gamma\left(1+\frac{n m-1}{p}\right)} \mathbb{E}[Z]=p n^{\frac{1}{p}} m^{\frac{1}{p}} \mathbb{E}[Z], \tag{381}
\end{equation*}
$$

where for (381) we introduce the following notation.

$$
\begin{equation*}
\mathrm{Z} \stackrel{\text { def }}{=}\left(\sum_{i=1}^{n} \mathrm{R}_{i} \mathrm{~N}_{i}\right)^{\frac{1}{2}} . \tag{382}
\end{equation*}
$$

Let R be a random variable that takes values in $[0, \infty)$ whose density at each $t \in(0, \infty)$ is given by (353), i.e., $R_{1}, \ldots, R_{n}$ are independent copies of $R$. We computed the moments of $R$ in (363) and by Stirling's formula this gives the following asymptotic evaluations.

$$
\begin{align*}
\mathbb{E}\left[\mathrm{R}^{\frac{1}{2}}\right] & =\frac{m^{1-\frac{1}{p}}}{p}  \tag{383}\\
\mathbb{E}[\mathrm{R}] & =\max \left\{\frac{m}{p}, 1\right\} \frac{m^{1-\frac{2}{p}}}{p}  \tag{384}\\
\mathbb{E}\left[\mathrm{R}^{2}\right] & =\max \left\{\frac{m^{3}}{p^{3}}, 1\right\} \frac{m^{1-\frac{4}{p}}}{p} . \tag{385}
\end{align*}
$$

We also need an analogous asymptotic evaluation of moments of the random variable N in (380). Observe that the random variables N and $\|\mathrm{U}\|_{\ell_{q}^{m}}$ are independent, since $\mathrm{U} /\|\mathrm{U}\|_{\ell_{q}^{m}}$ and $\|\mathrm{U}\|_{\ell_{q}^{m}}$ are independent and $N$ is a function $\mathrm{U} /\|\mathrm{U}\|_{\ell_{q}^{m}}$. Consequently, for every $\beta>0$ we have

$$
\begin{equation*}
\mathbb{E}\left[\|\mathrm{U}\|_{\ell_{q}^{m}}^{(2 q-2) \beta}\right] \mathbb{E}\left[\mathrm{N}^{\beta}\right]=\mathbb{E}\left[\|\mathrm{U}\|_{\ell_{q}^{m}}^{(2 q-2) \beta} \mathrm{N}^{\beta}\right] \stackrel{[380}{=} \mathbb{E}\left[\|\mathrm{U}\|_{\ell_{2 q-2}^{m}}^{(2 q-2) \beta}\right] . \tag{386}
\end{equation*}
$$

Since (e.g. by Lemma 155 the density of $\|U\|_{\ell_{q}^{m}}$ at $s \in(0, \infty)$ is proportional to $s^{m-1} e^{-s^{q}}$, we can compute analogously to (359) that

$$
\mathbb{E}\left[\|U\|_{\ell_{q}^{m}}^{(2 q-2) \beta}\right]=\frac{\int_{0}^{\infty} s^{m-1+(2 q-2) \beta} e^{-s^{q}} \mathrm{~d} s}{\int_{0}^{\infty} r^{m-1} e^{-r^{q}} \mathrm{~d} r}=\frac{\Gamma\left(2 \beta+\frac{m-2 \beta}{q}\right)}{\Gamma\left(\frac{m}{q}\right)} .
$$

Therefore (386) implies that

$$
\mathbb{E}\left[N^{\beta}\right]=\frac{\Gamma\left(\frac{m}{q}\right)}{\Gamma\left(2 \beta+\frac{m-2 \beta}{q}\right)} \mathbb{E}\left[\|U\|_{\ell_{2 q-2}}^{(2 q-2) \beta}\right] .
$$

By considering each of the values $\beta \in\left\{\frac{1}{2}, 1,2\right\}$ in this identity and using Stirling's formula, we get the following asymptotic evaluations of moments of N in terms of moments of $\|\mathrm{U}\|_{\ell_{2 q-2}^{m}}$.

$$
\begin{align*}
\mathbb{E}\left[N^{\frac{1}{2}}\right] & =\frac{q}{m^{1-\frac{1}{q}}} \mathbb{E}\left[\|U\|_{\ell_{2 q-2}^{m}}^{q-1}\right],  \tag{387}\\
\mathbb{E}[\mathrm{N}] & =\min \left\{\frac{q}{m}, 1\right\} \frac{q}{m^{1-\frac{2}{q}}} \mathbb{E}\left[\|U\|_{\ell_{2 q-2}^{m}}^{2 q-2}\right],  \tag{388}\\
\mathbb{E}\left[N^{2}\right] & =\min \left\{\frac{q^{3}}{m^{3}}, 1\right\} \frac{q}{m^{1-\frac{4}{q}}} \mathbb{E}\left[\|U\|_{\ell_{2 q-2}^{m}}^{4-4}\right] . \tag{389}
\end{align*}
$$

Due to (387), (388), (389), we will next evaluate the corresponding moments of $\|\mathrm{U}\|_{\ell_{2 q-2}^{m}}$. Recalling the density (379) of G, for every $\beta>-1 /(2 q-2)$ we have

$$
\mathbb{E}\left[|\mathrm{G}|^{(2 q-2) \beta}\right]=\frac{1}{\Gamma\left(1+\frac{1}{q}\right)} \int_{0}^{\infty} s^{(2 q-2) \beta} e^{-s^{q}} \mathrm{~d} s=\frac{\Gamma\left(\frac{2 q-2}{q} \beta+\frac{1}{q}\right)}{q \Gamma\left(1+\frac{1}{q}\right)} .
$$

Hence,

$$
\begin{equation*}
\mathbb{E}\left[|\mathrm{G}|^{q-1}\right]=\mathbb{E}\left[|\mathrm{G}|^{2 q-2}\right]=\mathbb{E}\left[|\mathrm{G}|^{4 q-4}\right]=\frac{1}{q} . \tag{390}
\end{equation*}
$$

We therefore have

$$
\begin{equation*}
\mathbb{E}\left[\|\mathrm{U}\|_{\ell_{2 q-2}^{m}}^{2 q-2}\right]=m \mathbb{E}\left[|\mathrm{G}|^{2 q-2}\right] \stackrel{390}{=} \frac{m}{q}, \tag{391}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathbb{E}\left[\|\mathrm{U}\|_{\ell_{2 q-2}}^{4 q-4}\right]=\mathbb{E}\left[\left(\sum_{j=1}^{m}\left|\mathrm{G}_{j}\right|^{2 q-2}\right)^{2}\right]=m \mathbb{E}\left[|\mathrm{G}|^{4 q-4}\right]+m(m-1)\left(\mathbb{E}\left[|G|^{2 q-2}\right]\right)^{2} \stackrel{[390}{-} \max \left\{\frac{m}{q}, 1\right\} \frac{m}{q} . \tag{392}
\end{equation*}
$$

Consequently, using Hölder's inequality we get the following estimate.

$$
\begin{align*}
& \frac{m}{q} \stackrel{(391)}{=} \mathbb{E}\left[\|U\|_{\ell_{2 q-2}^{m}}^{2 q-2}\right]=\mathbb{E}\left[\|U\|_{\ell_{2 q-2}^{m}}^{\frac{2}{3}(q-1)}\|U\|_{\ell_{2 q-2}^{m}}^{\frac{1}{3}(4 q-4)}\right] \\
& \leqslant\left(\mathbb{E}\left[\|U\|_{\ell_{2 q-2}^{m}}^{q-1}\right]\right)^{\frac{2}{3}}\left(\mathbb{E}\left[\|U\|_{\ell_{2 q-2}^{m}}^{4 q-4}\right]\right)^{\frac{1}{3}} \stackrel{(392)}{=}\left(\mathbb{E}\left[\|U\|_{\ell_{2 q-2}^{m}}^{q-1}\right]\right)^{\frac{2}{3}}\left(\max \left\{\frac{m}{q}, 1\right\} \frac{m}{q}\right)^{\frac{1}{3}} . \tag{393}
\end{align*}
$$

This simplifies to give

$$
\begin{equation*}
\mathbb{E}\left[\|\cup\|_{\ell_{2 q-2}^{m}}^{q-1}\right] \gtrsim \min \left\{\sqrt{\frac{m}{q}}, \frac{m}{q}\right\} . \tag{394}
\end{equation*}
$$

At the same time, by Cauchy-Schwarz,

$$
\begin{equation*}
\mathbb{E}\left[\|\cup\|_{\ell_{2 q-2}^{m}}^{q-1}\right] \leqslant\left(\mathbb{E}\left[\|U\|_{\ell_{2 q-2}^{m}}^{2 q-2}\right]\right)^{\frac{1}{2}} \stackrel{(391)}{=} \sqrt{\frac{m}{q}} . \tag{395}
\end{equation*}
$$

Also, by the subadditivity of the square root on $[0, \infty)$,

$$
\begin{equation*}
\mathbb{E}\left[\|\cup\|_{\ell_{2 q-2}}^{q-1}\right]=\mathbb{E}\left[\left(\sum_{j=1}^{m}\left|\mathrm{G}_{j}\right|^{2 q-2}\right)^{\frac{1}{2}}\right] \leqslant \mathbb{E}\left[\sum_{j=1}^{m}\left|\mathrm{G}_{j}\right|^{q-1}\right]=m \mathbb{E}\left[|\mathrm{G}|^{q-1}\right] \stackrel{\sqrt{390}}{=} \frac{m}{q} . \tag{396}
\end{equation*}
$$

By combining (395) and (396) we see that (394) is in fact sharp, i.e.,

$$
\begin{equation*}
\mathbb{E}\left[\|\cup\|_{\ell_{2 q-2}^{m}}^{q-1}\right]=\min \left\{\sqrt{\frac{m}{q}}, \frac{m}{q}\right\} . \tag{397}
\end{equation*}
$$

By substituting (397) into (387), and correspondingly (391) into (388) and (392) into (389), we get the following asymptotic identities.

$$
\begin{align*}
\mathbb{E}\left[\mathrm{N}^{\frac{1}{2}}\right] & =\min \left\{\sqrt{\frac{q}{m}}, 1\right\} m^{\frac{1}{q}},  \tag{398}\\
\mathbb{E}[\mathrm{~N}] & =\min \left\{\frac{q}{m}, 1\right\} m^{\frac{2}{q}},  \tag{399}\\
\mathbb{E}\left[\mathrm{~N}^{2}\right] & =\min \left\{\frac{q^{2}}{m^{2}}, 1\right\} m^{\frac{4}{q}} \tag{400}
\end{align*}
$$

By combining (384) and (399) we see that

$$
\mathbb{E}\left[Z^{2}\right]=n(\mathbb{E}[R])(\mathbb{E}[\mathrm{N}])=\frac{\max \{m, p\} \min \{q, m\}}{p^{2}} n m^{\frac{2}{q}-\frac{2}{p}} .
$$

Using Cauchy-Schwarz, this implies the following upper bound on the final term in (381).

$$
\begin{equation*}
p n^{\frac{1}{p}} m^{\frac{1}{p}} \mathbb{E}[Z] \leqslant p n^{\frac{1}{p}} m^{\frac{1}{p}}\left(\mathbb{E}\left[\mathrm{Z}^{2}\right]\right)^{\frac{1}{2}}=n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{q}} \sqrt{\max \{m, p\} \min \{m, q\}} . \tag{401}
\end{equation*}
$$

Also, recalling (382) and using the subadditivity of the square root on $[0, \infty$ ) in combination with (383) and (398), we have the following additional upper bound on the final term in (381).

$$
\begin{equation*}
p n^{\frac{1}{p}} m^{\frac{1}{p}} \mathbb{E}[Z] \leqslant p n^{\frac{1}{p}} m^{\frac{1}{p}} \mathbb{E}\left[\sum_{i=1}^{n} \mathrm{R}_{i}^{\frac{1}{2}} \mathrm{~N}_{i}^{\frac{1}{2}}\right]=p n^{1+\frac{1}{p}} m^{\frac{1}{p}}\left(\mathbb{E}\left[\mathrm{R}^{\frac{1}{2}}\right]\right)\left(\mathbb{E}\left[\mathrm{N}^{\frac{1}{2}}\right]\right)=n^{1+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} \sqrt{\min \{m, q\}} . \tag{402}
\end{equation*}
$$

It follows from (401) and (402) that

$$
\begin{align*}
p n^{\frac{1}{p}} m^{\frac{1}{p}} \mathbb{E}[Z] & \lesssim n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{q}} \sqrt{\min \{m, q\}} \min \{\sqrt{n m}, \sqrt{\max \{m, p\}}\} \\
& = \begin{cases}n^{1+\frac{1}{p}} m^{1+\frac{1}{q}} & m \leqslant \min \left\{\frac{p}{n}, q\right\}, \\
\sqrt{q} n^{1+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} & q \leqslant m \leqslant \frac{p}{n}, \\
\sqrt{p} n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} & \frac{p}{n} \leqslant m \leqslant \min \{p, q\}, \\
\sqrt{p q} n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{q}} & \max \left\{\frac{p}{n}, q\right\} \leqslant m \leqslant p, \\
n^{\frac{1}{2}+\frac{1}{p}} m^{1+\frac{1}{q}} & p \leqslant m \leqslant q, \\
\sqrt{q} n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} & m \geqslant \max \{p, q\} .\end{cases} \tag{403}
\end{align*}
$$

We will next prove that (403) is optimal in all of the six ranges that appear in (403); by (381) and (316), this will complete the proof of Corollary 160 Recalling (382) and using (384, (385), (399), (400), the fourth moment of $Z$ can be evaluated (up to universal constant factors) as follows.

$$
\begin{align*}
\mathbb{E}\left[Z^{4}\right] & =\mathbb{E}\left[\sum_{i=1}^{n} \sum_{j=1}^{n} \mathrm{R}_{i} \mathrm{R}_{j} \mathrm{~N}_{i} \mathrm{~N}_{j}\right] \\
& =n\left(\mathbb{E}\left[\mathrm{R}^{2}\right]\right)\left(\mathbb{E}\left[\mathrm{N}^{2}\right]\right)+n(n-1)(\mathbb{E}[\mathrm{R}])^{2}(\mathbb{E}[\mathrm{~N}])^{2} \\
& =\frac{(\max \{m, p\})^{3}(\min \{m, q\})^{2}}{p^{4}} n m^{\frac{4}{q}-\frac{4}{p}-1}+\frac{(\max \{m, p\} \min \{m, q\})^{2}}{p^{4}} n^{2} m^{\frac{4}{q}-\frac{4}{p}}  \tag{404}\\
& =\frac{(\max \{m, p\} \min \{m, q\})^{2} \max \{n m, p\}}{p^{4}} n m^{\frac{4}{q}-\frac{4}{p}-1} .
\end{align*}
$$

By using Hölder's inequality similarly to (393), we conclude that

$$
\begin{aligned}
p n^{\frac{1}{p}} m^{\frac{1}{p}} \mathbb{E}[Z] \geqslant p n^{\frac{1}{p}} m^{\frac{1}{p}} \frac{\left(\mathbb{E}\left[\mathrm{Z}^{2}\right]\right)^{\frac{3}{2}}}{\left(\mathbb{E}\left[\mathrm{Z}^{4}\right]\right)^{\frac{1}{2}}} \\
\stackrel{401)}{\wedge} \xlongequal{404} n^{1+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} \frac{\sqrt{\max \{m, p\} \min \{m, q\}}}{\sqrt{\max \{n m, p\}}}= \begin{cases}n^{1+\frac{1}{p}} m^{1+\frac{1}{q}} & m \leqslant \min \left\{\frac{p}{n}, q\right\}, \\
\sqrt{q} n^{1+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} & q \leqslant m \leqslant \frac{p}{n}, \\
\sqrt{p} n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} & \frac{p}{n} \leqslant m \leqslant \min \{p, q\}, \\
\sqrt{p q} n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{q}} & \max \left\{\frac{p}{n}, q\right\} \leqslant m \leqslant p, \\
n^{\frac{1}{2}+\frac{1}{p}} m^{1+\frac{1}{q}} & p \leqslant m \leqslant q, \\
\sqrt{q} n^{\frac{1}{2}+\frac{1}{p}} m^{\frac{1}{2}+\frac{1}{q}} & m \geqslant \max \{p, q\} .\end{cases}
\end{aligned}
$$

Lemma 161 below applies Theorem 160 iteratively to obtain an upper bound on the surface area of the unit sphere of nested $\ell_{p}$ norms on $k$-tensors (the case $k=2$ corresponds to $n$ by $m$ matrices equipped with the $\ell_{p}^{n}\left(\ell_{q}^{m}\right)$ norm). The second part of Lemma 161 namely the conclusion (406) below, is an implementation of the approach towards Conjecture 9 for the hypercube that we described in Remark 55 .

Lemma 161. Suppose that $k, n_{1}, \ldots, n_{k} \in \mathbb{N}$ and $p_{1}, \ldots, p_{k} \in[1, \infty]$ are such that $n_{1} \geqslant \max \left\{3, p_{1}-2\right\}$ and $n_{1} n_{2} \ldots n_{j-1} \geqslant p_{j}-2$ for every $j \in\{2, \ldots, k\}$. Define normed spaces $\mathbf{Y}_{0}, \mathbf{Y}_{1}, \ldots, \mathbf{Y}_{k}$ by setting $\mathbf{Y}_{0}=\mathbb{R}$ and inductively $\mathbf{Y}_{j}=\ell_{p_{j}}^{n_{j}}\left(\mathbf{Y}_{j-1}\right)$ for $j \in\{1, \ldots, k\}$. Then,

$$
\begin{equation*}
\frac{\operatorname{vol}_{n_{1} \ldots n_{k}-1}\left(\partial B_{\mathbf{Y}_{k}}\right)}{\operatorname{vol}_{n_{1} \ldots n_{k}}\left(B_{\mathbf{Y}_{k}}\right)} \leqslant e^{O(k)} \sqrt{p_{1}} \prod_{j=1}^{k} n_{j}^{\frac{1}{2}+\frac{1}{p_{j}}} . \tag{405}
\end{equation*}
$$

Hence, using the natural identification of the vector space that underlies $\mathbf{Y}_{k}$ with $\mathbb{R}^{\operatorname{dim}\left(\mathbf{Y}_{k}\right)}=\mathbb{R}^{n_{1} n_{2} \ldots n_{k}}$, if in addition we have $n_{1}=O(1)$ and $p_{j}=\log n_{j}$ for every $j \in\{1, \ldots, k\}$, then

$$
\begin{equation*}
B_{\mathbf{Y}_{k}} \subseteq B_{\ell_{\infty}^{\operatorname{dim}\left(\mathbf{Y}_{k}\right)}} \subseteq e^{O(k)} B_{\mathbf{Y}_{k}} \quad \text { and } \quad \frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}_{k}}\right)}{\operatorname{vol}_{\operatorname{dim}\left(\mathbf{Y}_{k}\right)}\left(B_{\mathbf{Y}_{k}}\right)} \leqslant e^{O(k)} \tag{406}
\end{equation*}
$$

where we recall the notation (52).
Proof. Suppose that $n, m \in \mathbb{N}$ and $p \in(1, \infty)$. By applying Cauchy-Schwarz to the right hand side of (354) while using the case $\alpha=1$ of (363), we see that for every normed space $\mathbf{X}=\left(\mathbb{R}^{m},\|\cdot\|_{\mathbf{X}}\right)$ we have

$$
\begin{equation*}
\frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}(\mathbf{X})}\right)}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}(\mathbf{X})}\right)} \leqslant \frac{p \Gamma\left(1+\frac{n m}{p}\right)}{\Gamma\left(1+\frac{n m-1}{p}\right)}\left(\frac{n \Gamma\left(\frac{m+2 p-2}{p}\right)}{\Gamma\left(\frac{m}{p}\right)} \int_{\partial B_{\mathbf{X}}}\|\nabla\| \cdot\left\|_{\mathbf{X}}\right\|_{\ell_{2}^{m}}^{2} \mathrm{~d} \kappa_{\mathbf{X}}\right)^{\frac{1}{2}} \tag{407}
\end{equation*}
$$

If also $m \geqslant \max \{3, p-2\}$, then by Stirling's formula 407) gives the following estimate.

$$
\begin{equation*}
\frac{\operatorname{vol}_{n m-1}\left(\partial B_{\ell_{p}^{n}(\mathbf{X})}\right)}{\operatorname{vol}_{n m}\left(B_{\ell_{p}^{n}(\mathbf{X})}\right)} \lesssim n^{\frac{1}{2}+\frac{1}{p}} m\left(\int_{\partial B_{\mathbf{X}}}\|\nabla\| \cdot\left\|_{\mathbf{X}}\right\|_{\ell_{2}^{m}}^{2} \mathrm{~d} \kappa_{\mathbf{X}}\right)^{\frac{1}{2}} \tag{408}
\end{equation*}
$$

By continuity we may assume that $p_{1}, \ldots, p_{k} \in(1, \infty)$. Denote $d_{0}=1$ and $d_{j}=\operatorname{dim}\left(\mathbf{Y}_{j}\right)=n_{1} n_{2} \ldots n_{j}$ for $j \in\{1, \ldots, k\}$. We will naturally identify $\mathbf{Y}_{j}$ with $\left(\mathbb{R}^{d_{j}},\|\cdot\|_{\mathbf{Y}_{j}}\right)$. As $\mathbf{Y}_{k}=\ell_{p_{k}}^{n_{k}}\left(\mathbf{Y}_{k-1}\right)$, we deduce from 408) that

$$
\begin{equation*}
\frac{\operatorname{vol}_{d_{k}-1}\left(\partial B_{\mathbf{Y}_{k}}\right)}{\operatorname{vol}_{d_{k}}\left(B_{\mathbf{Y}_{k}}\right)} \lesssim n_{k}^{\frac{1}{2}+\frac{1}{p_{k}}}\left(\prod_{j=1}^{k-1} n_{j}\right)\left(\int_{\partial B_{\mathbf{Y}_{k-1}}}\|\nabla\| \cdot\left\|_{\mathbf{Y}_{k-1}}\right\|_{\ell_{2}^{d_{k-1}}}^{2} \mathrm{~d} \kappa_{\mathbf{Y}_{k-1}}\right)^{\frac{1}{2}} \tag{409}
\end{equation*}
$$

At the same time, by (355) for every $j \in\{1, \ldots, k\}$ we have

$$
\begin{equation*}
\int_{\partial B_{\mathbf{Y}_{j}}}\|\nabla\| \cdot\left\|_{\mathbf{Y}_{j}}\right\|_{\ell_{2}^{d_{j}}}^{2} \mathrm{~d} \kappa_{\mathbf{Y}_{j}}=\frac{n_{j} \Gamma\left(\frac{d_{j}}{p_{j}}\right) \Gamma\left(\frac{d_{j-1}+2 p_{j}-2}{p_{j}}\right)}{\Gamma\left(\frac{d_{j-1}}{p_{j}}\right) \Gamma\left(\frac{d_{j}+2 p_{j}-2}{p_{j}}\right)} \int_{\partial B_{\mathbf{Y}_{j-1}}}\|\nabla\| \cdot\left\|_{\mathbf{Y}_{j-1}}\right\|_{\ell_{2}^{d_{j-1}}}^{2} \mathrm{~d} \kappa_{\mathbf{Y}_{j-1}} \tag{410}
\end{equation*}
$$

If also $j \geqslant 2$, then $d_{j-1} \geqslant n_{1} \geqslant 3$ and by assumption $d_{j-1} \geqslant p_{j}-2$, so by Stirling's formula 410) gives

$$
\begin{equation*}
\forall j \in\{2, \ldots, k\}, \quad \int_{\partial B_{\mathbf{Y}_{j}}}\|\nabla\| \cdot\left\|_{\mathbf{Y}_{j}}\right\|_{\ell_{2}^{d_{j}}}^{2} \mathrm{~d} \kappa_{\mathbf{Y}_{j}}=n_{j}^{\frac{2}{p_{j}}-1} \int_{\partial B_{\mathbf{Y}_{j-1}}}\|\nabla\| \cdot\left\|_{\mathbf{Y}_{j-1}}\right\|_{\ell_{2}^{d_{j-1}}}^{2} \mathrm{~d} \kappa_{\mathbf{Y}_{j-1}} \tag{411}
\end{equation*}
$$

When $j=1$ we have $d_{0}=1$ and $n_{1} \geqslant \max \left\{3, p_{1}-2\right\}$, and therefore by Stirling's formula (410) gives

$$
\begin{equation*}
\int_{\partial B_{\mathbf{Y}_{1}}}\|\nabla\| \cdot\left\|_{\mathbf{Y}_{1}}\right\|_{\ell_{2}^{d_{1}}}^{2} \mathrm{~d} \kappa_{\mathbf{Y}_{1}}=p_{1} n_{1}^{\frac{2}{p_{1}}-1} \tag{412}
\end{equation*}
$$

Hence, by applying (411) iteratively in combination with the base case 412), we conclude that

$$
\begin{equation*}
\int_{\partial B_{\mathbf{Y}_{k-1}}}\|\nabla\| \cdot\left\|_{\mathbf{Y}_{k-1}}\right\|_{\ell_{2}^{d_{k-1}}}^{2} \mathrm{~d} \kappa_{\mathbf{Y}_{k-1}} \leqslant e^{O(k)} p_{1} \prod_{j=1}^{k-1} n_{j}^{\frac{2}{p_{j}}-1} \tag{413}
\end{equation*}
$$

A substitution of (413) into (409) yields the desired estimate 405).
To deduce the conclusion (406), note that for every $j \in\{1, \ldots, k\}$ we have the point-wise bounds

$$
\|\cdot\|_{\ell_{\infty}^{n_{j}}\left(\mathbf{Y}_{j-1}\right)} \leqslant\|\cdot\|_{\mathbf{Y}_{j}}=\|\cdot\|_{\ell_{p_{j}}^{n_{j}}\left(\mathbf{Y}_{j-1}\right)} \leqslant n_{j}^{\frac{1}{p_{j}}}\|\cdot\|_{\ell_{\infty}^{n_{j}}\left(\mathbf{Y}_{j-1}\right)} .
$$

It follows by induction that

$$
\|\cdot\|_{\ell_{\infty}^{d_{k}}} \leqslant\|\cdot\|_{\mathbf{Y}_{k}} \leqslant\left(\prod_{j=1}^{k} n_{j}^{\frac{1}{p_{j}}}\right)\|\cdot\|_{\ell_{\infty}^{d_{k}}}=e^{O(k)}\|\cdot\|_{\ell_{\infty}^{d_{k}}},
$$

where the final step holds if $p_{j}=\log n_{j}$ for every $j \in\{1, \ldots, k\}$. This implies the inclusions in 406). Furthermore, $\mathbf{Y}_{k}$ belongs to the class of spaces from Example 39. Hence $\mathbf{Y}_{k}$ is canonically positioned and by the discussion in Section 1.6 .2 know that $B_{\mathbf{Y}^{\prime}}$ is in its minimum surface area position. Therefore,

$$
\frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}_{k}}\right)}{\operatorname{vol}_{d_{k}}\left(B_{\mathbf{Y}_{k}}\right)}=\frac{\operatorname{vol}_{d_{k}-1}\left(\partial B_{\mathbf{Y}_{k}}\right)}{\operatorname{vol}_{d_{k}}\left(B_{\mathbf{Y}_{k}}\right) \sqrt{d_{k}}} \leqslant e^{O(k)} \sqrt{p_{1}} \prod_{j=1}^{k} n_{j}^{\frac{1}{p_{j}}}=e^{O(k)},
$$

where the first step uses [GP99, Proposition 3.1], the second step is 405], and the final step holds because $p_{1}=O(1)$ and $p_{j}=\log n_{j}$. This completes the proof of 406).

The following technical lemma replaces a more ad-hoc argument that we previously had to deduce Proposition 163 below from Lemma 161 it is due to Noga Alon and we thank him for allowing us to include it here. This lemma shows that the set of super-lacunary products $n_{1} n_{2} \ldots n_{k}$ that can serve as dimensions of the space $\mathbf{Y}_{k}$ in Lemma 161 for which (406) holds is quite dense in $\mathbb{N}$.

Lemma 162. For every integer $n \geqslant 3$ there are $k, m \in \mathbb{N} \cup\{0\}$ and integers $n_{1}<n_{2}<\ldots<n_{k}$ that satisfy

- $n=n_{1} n_{2} \ldots n_{k}+m$,
- $n_{1} \in\{6,7\}$ and $n_{i+1} \leqslant 2^{n_{i}} \leqslant n_{i+1}^{3}$ for every $i \in\{1, \ldots, k-1\}$,
- $m \leqslant(\log n)^{1+o(1)}$.

Prior to proving Lemma 162, we will make some preparatory (mechanical) observations for ease of later reference. Note first that the conclusion $n_{i+1} \leqslant 2^{n_{i}} \leqslant n_{i+1}^{3}$ of Lemma 162 can be rewritten as

$$
\forall i \in\{1, \ldots, k-1\}, \quad \log _{2} n_{i+1} \leqslant n_{i} \leqslant \log _{\sqrt[3]{2}} n_{i+1} .
$$

It follows by induction that

$$
\begin{equation*}
\forall i \in\{1, \ldots, k\}, \quad \log _{2}^{[k-i]} n_{k} \leqslant n_{i} \leqslant \log _{\sqrt[3]{2}}^{[k-i]} n_{k}, \tag{414}
\end{equation*}
$$

where, as in (135], we denote the iterates of $\varphi:(0, \infty) \rightarrow \mathbb{R}$ by $\varphi^{[j]}=\varphi \circ \varphi^{[j-1]}:\left(\varphi^{[j-1]}\right)^{-1}(0, \infty) \rightarrow \mathbb{R}$ for each $j \in \mathbb{N}$, with the convention $\varphi^{[0]}(x)=x$ for every $x \in(0, \infty)$. Since $n_{1} \in\{6,7\}$, it follows from (414) that

$$
\begin{equation*}
k=\log ^{*} n_{k} \lesssim \log ^{*} n . \tag{415}
\end{equation*}
$$

Consequently,

$$
\begin{aligned}
n_{k} \log n_{k}=n_{k} n_{k-1} \leqslant \prod_{i=1}^{k} n_{k} \leqslant n=m & +\prod_{i=1}^{k} n_{k} \leqslant(\log n)^{1+o(1)}+\prod_{i=1}^{k} \log _{\sqrt[3]{2}}^{[k-i]} n_{k} \\
& \lesssim(\log n)^{2}+n_{k}\left(\log n_{k}\right)\left(\log \log n_{k}\right)^{O\left(\log ^{*} n_{k}\right)} \lesssim(\log n)^{2}+n_{k}\left(\log n_{k}\right)^{2} .
\end{aligned}
$$

This implies the following (quite crude) bounds on $n_{k}$.

$$
\begin{equation*}
\frac{n}{(\log n)^{2}} \lesssim n_{k} \lesssim \frac{n}{\log n} . \tag{416}
\end{equation*}
$$

Note in particular that thanks to (416) we know that (415) can be improved to $k=\log ^{*} n$.
Proof of Lemma 162 . Let $\mathbb{M} \subseteq \mathbb{N}$ be the set of all those $x \in \mathbb{N}$ that can be written as $x=n_{1} n_{2} \ldots n_{k}$ for some $k, n_{1}, \ldots, n_{k} \in \mathbb{N}$ that satisfy $n_{k}>n_{k-1}>\ldots>n_{1} \in\{6,7\}$ and

$$
\begin{equation*}
\forall i \in\{1, \ldots, k-1\}, \quad n_{i+1} \leqslant 2^{n_{i}} \leqslant n_{i+1}^{3} . \tag{417}
\end{equation*}
$$

The goal of Lemma 162 is to show that there exists $x \in \mathbb{M}$ such that

$$
\begin{equation*}
n-(\log n)^{1+o(1)} \leqslant x \leqslant n . \tag{418}
\end{equation*}
$$

By adjusting the $o(1)$ term, we may assume that $n$ is sufficiently large, say, $n \geqslant n(0)$ for some fixed $n(0) \in \mathbb{N}$ that will be determined later. We will then find $x \in \mathbb{M}$ with a representation $x=n_{1} n_{2} \ldots n_{k}$ as above and

$$
\begin{equation*}
n-n_{1} n_{2} \ldots n_{k-1} \leqslant x \leqslant n . \tag{419}
\end{equation*}
$$

This would imply the desired bound 418 because

$$
\begin{equation*}
\prod_{i=1}^{k-1} n_{i} \stackrel{\sqrt[414]{\gtrless}}{\lessgtr} \prod_{i=1}^{k-1} \log _{\sqrt[3]{2}}^{[k-i]} n_{k} \stackrel{\sqrt[415]{\lesssim}}{\lesssim}\left(\log n_{k}\right)^{1+o(1)} \stackrel{\sqrt[416]{\lesssim}}{\lesssim}(\log n)^{1+o(1)} \tag{420}
\end{equation*}
$$

We will first construct $\left\{y_{i}\right\}_{i=1}^{\infty} \subseteq \mathbb{M}$ such that $y_{1}=7$ and $y_{i}<y_{i+1}<12 y_{i}$ for every $i \in \mathbb{N}$. Furthermore, for each $i \in \mathbb{N}$ there are $k, n_{1}, \ldots, n_{k} \in \mathbb{N}$ with $y_{i}=n_{1} n_{2} \ldots n_{k}$ such that $n_{k}>n_{k-1}>\ldots>n_{1} \in\{6,7\}$ and

$$
\begin{equation*}
\forall j \in\{1, \ldots, k-1\}, \quad n_{j+1}^{2} \leqslant 2^{n_{j}} \leqslant 2 n_{j+1}^{2}, \tag{421}
\end{equation*}
$$

which is a more stringent requirement than 417). Note in passing that 421) implies the (crude) bound

$$
\begin{equation*}
\prod_{j=1}^{k}\left(1+\frac{1}{n_{j}}\right) \leqslant 2 \tag{422}
\end{equation*}
$$

To verify (422), note that since $\left\{n_{j}\right\}_{j=1}^{k}$ is strictly increasing and the second inequality in (421) holds, it is mechanical to check that $n_{1} \geqslant 6, n_{2} \geqslant 7, n_{3} \geqslant 8, n_{4} \geqslant 12$ and $n_{j+1} \geqslant 3 n_{j}$ for every $j \in\{4,5, \ldots, k-1\}$. So,

$$
\prod_{j=1}^{k}\left(1+\frac{1}{n_{j}}\right) \leqslant\left(1+\frac{1}{6}\right)\left(1+\frac{1}{7}\right)\left(1+\frac{1}{8}\right) e^{\sum_{s=0}^{\infty} \frac{1}{12 \cdot 3^{s}}}=\left(1+\frac{1}{6}\right)\left(1+\frac{1}{7}\right)\left(1+\frac{1}{8}\right) e^{\frac{1}{8}} \leqslant 2
$$

Suppose that $y_{i}$ has been defined with a representation $y_{i}=n_{1} n_{2} \ldots n_{k}$ that fulfils the above requirements. Define $m_{0}, m_{1}, \ldots, m_{k} \in \mathbb{N}$ with $m_{0}=6, m_{k}=n_{k}+1$ and $m_{j} \in\left\{n_{j}, n_{j}+1\right\}$ for all $j \in\{1, \ldots, k-1\}$ by induction as follows. Assuming that $m_{j+1}$ has already been constructed for some $j \in\{1, \ldots, k-1\}$, let

$$
m_{j} \stackrel{\text { def }}{=} \begin{cases}n_{j} & \text { if } m_{j+1}^{2} \leqslant 2^{n_{j}}  \tag{423}\\ n_{j}+1 & \text { if } m_{j+1}^{2}>2^{n_{j}}\end{cases}
$$

Definition 423) implies that $m_{j}<m_{j+1}$. Indeed, $n_{j}<n_{j+1}$ so if $m_{j}=n_{j}$, then $n_{j}<n_{j+1} \leqslant m_{j+1}$ since $m_{j+1} \geqslant n_{j+1}$ by the induction hypothesis. On the other hand, if $m_{j}=n_{j}+1$, then since the first inequality in (421) holds, the definition (423) necessitates that $m_{j+1}=n_{j}+1$, so $m_{j}<m_{j+1}$ in this case as well.

Definition (423) also ensures that the requirement (421) is inherited by $\left\{m_{j}\right\}_{j=1}^{k}$, i.e.,

$$
\begin{equation*}
\forall j \in\{1, \ldots, k-1\}, \quad m_{j+1}^{2} \leqslant 2^{m_{j}} \leqslant 2 m_{j+1}^{2} \tag{424}
\end{equation*}
$$

Indeed, if $m_{j}=n_{j}$, then $m_{j+1}^{2} \leqslant 2^{n_{j}}=2^{m_{j}}$ by 423, i.e., the first inequality in (424) holds, and the second inequality in (424) holds because $m_{j+1} \geqslant n_{j+1}$ and (421) holds. On the other hand, if $m_{j}=n_{j}+1$, then by (423) we necessarily have $m_{j+1}=n_{j}+1$ and $m_{j+1}^{2}>2^{n_{j}}$, which directly gives the second inequality in (424), and in combination with (421) we also get the first inequality in (424) because

$$
\frac{m_{j+1}}{2^{m_{j}}}=\frac{\left(n_{j}+1\right)^{2}}{2^{n_{j}+1}} \stackrel{421}{\leqslant} \frac{\left(n_{j}+1\right)^{2}}{2 n_{j}^{2}} \leqslant 1
$$

where the final step uses $n_{j} \geqslant 6$, though $n_{j} \geqslant 1 /(\sqrt{2}-1)=2.414 \ldots$ is all that is needed for this purpose.
If the above construction produces $m_{1} \in\{6,7\}$, then define $y_{i+1}=m_{1} m_{2} \ldots m_{k}$. Otherwise necessarily $m_{1}=n_{1}+1=8$, so (424) holds also when $j=0$ (recall that $m_{0}=6$, hence $m_{1}^{2}=2^{6}=2^{m_{0}}$ ), so we can define $y_{i+1}=m_{0} m_{1} \ldots m_{k}$ and thanks to (424) in both cases $y_{i+1}$ has the desired form. Moreover,

$$
\frac{y_{i+1}}{y_{i}} \leqslant 6 \prod_{j=1}^{k}\left(1+\frac{1}{n_{j}}\right) \stackrel{422}{\leqslant} 12 .
$$

This completes the inductive construction of the desired sequence $\left\{y_{i}\right\}_{i=1}^{\infty} \subseteq \mathbb{M}$.
With the sequence $\left\{y_{i}\right\}_{i=1}^{\infty} \subseteq \mathbb{M}$ at hand, will next explain how to obtain for each integer $n \geqslant n(0)$, where $n(0) \in \mathbb{N}$ is a sufficiently large universal constant that is yet to be determined, an element $x \in \mathbb{M}$ that approximates $n$ as in 419). Let $i \in \mathbb{N}$ be such that $y_{i} \leqslant n \leqslant y_{i+1}$ and denote $y=y_{i}$. Thus, there are $k, n_{1}, \ldots, n_{k} \in \mathbb{N}$ for which $y=n_{1} n_{2} \ldots n_{k}$ such that $n_{k}>n_{k-1}>\ldots>n_{1} \in\{6,7\}$ and (421) holds.

If $y \geqslant n-n_{1} n_{2} \ldots n_{k-1}$, then $x=y$ has the desired approximation property, so suppose from now that $y<n-n_{1} n_{2} \ldots n_{k-1}$, or equivalently $n /\left(n_{1} n_{2} \ldots n_{k-1}\right)>y /\left(n_{1} n_{2} \ldots n_{k-1}\right)+1=n_{k}+1$. Hence, if we define

$$
\begin{equation*}
n_{k}^{\prime} \stackrel{\text { def }}{=}\left\lfloor\frac{n}{n_{1} n_{2} \ldots, n_{k-1}}\right\rfloor \quad \text { and } \quad x=n_{1} n_{2} \ldots n_{k-1} n_{k}^{\prime} \tag{425}
\end{equation*}
$$

Then $n_{k}^{\prime} \geqslant n_{k}+1 \gtrsim n /(\log n)^{2}$, where we used (416). Consequently, recalling (414), there is a universal constant $n(0) \in \mathbb{N}$ such that if $n \geqslant n(0)$, then $n_{k}^{\prime}>\max \left\{144, n_{k-1}\right\}$. So, the sequence $n_{1}, n_{2}, \ldots, n_{k-1}, n_{k}^{\prime}$ is still increasing. Since by design $x$ satisfies 419), it remains to check that $x \in \mathbb{M}$, i.e., that 417) holds. Since $n_{1}, \ldots, n_{k}$ are assumed to satisfy the more stringent requirement 421, we only need to check that

$$
\begin{equation*}
n_{k}^{\prime} \leqslant 2^{n_{k-1}} \leqslant\left(n_{k}^{\prime}\right)^{3} \tag{426}
\end{equation*}
$$

The second inequality in (426) is valid since (421) holds and $n_{k}^{\prime}>n_{k}$. For the first inequality in (426), note that $y \leqslant n \leqslant 12 y$, as $y_{i+1} \leqslant 12 y_{i}$. Hence, $n_{k}^{\prime} \leqslant n /\left(n_{1} n_{2} \ldots n_{k-1}\right) \leqslant 12 y /\left(n_{1} n_{2} \ldots n_{k-1}\right)=12 n_{k}$. Therefore,

$$
2^{n_{k-1}} \stackrel{421}{\geqslant} n_{k}^{2} \geqslant\left(\frac{n_{k}^{\prime}}{12}\right)^{2}>n_{k}^{\prime},
$$

where the last step uses the fact that $n_{k}^{\prime}>144$.
We are now ready to extend the conclusion (406) of Lemma 161 to all dimensions $n \in \mathbb{N}$. Namely, we will prove the following proposition, which comes very close to proving Conjecture 9 for the hypercube $[-1,1]^{n}$ via a route that differs from the way by which we proved Theorem 24 .

Proposition 163. For any $n \in \mathbb{N}$ there is a normed space $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ that satisfies

$$
\forall x \in \mathbb{R}^{n} \backslash\{0\}, \quad\|x\|_{\ell_{\infty}^{n}} \leqslant\|x\|_{\mathbf{Y}} \leqslant e^{O\left(\log ^{*} n\right)}\|x\|_{\ell_{\infty}^{n}} \quad \text { and } \quad \frac{\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}} B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \leqslant e^{O\left(\log ^{*} n\right)}
$$

Furthermore, $\mathbf{Y}$ can be taken to be an $\ell_{\infty}$ direct sum of nested $\ell_{p}$ spaces as in Lemma 161 .
Proof. Let $\mathbb{M} \subseteq \mathbb{N}$ be the set of integers from the proof of Lemma 162 , namely $m \in \mathbb{M}$ if and only if there are integers $n_{k}>n_{k-1}>\ldots>n_{1} \in\{6,7\}$ that satisfy 417) such that $m=n_{1} n_{2} \ldots n_{k}$. By Lemma 161 , there exists $C>1$ such that for every $m \in \mathbb{M}$ there is a normed space $\mathbf{Y}^{m}=\left(\mathbb{R}^{m},\|\cdot\|_{\mathbf{Y}^{m}}\right)$ that satisfies

$$
\|\cdot\|_{\ell_{\infty}^{m}} \leqslant\|\cdot\|_{\mathbf{Y}^{m}} \leqslant e^{C \log ^{*} m}\|\cdot\|_{\ell_{\infty}^{m}} \quad \text { and } \quad \frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}^{m}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}^{m}}\right)} \leqslant e^{C \log ^{*} m}
$$

By applying Lemma 162 iteratively write $n=m_{1}+\ldots+m_{s+1}$ for $m_{1}, \ldots, m_{s} \in \mathbb{M}$ and $m_{s+1} \in\{1,2\}$ that satisfy $m_{i+1} \leqslant\left(\log m_{i}\right)^{c}$ for every $i \in\{1, \ldots, s\}$, where $c>1$ is a universal constant. Denote $\mathbf{Y}^{m_{s+1}}=\ell_{\infty}^{m_{s+1}}$ and consider the $\ell_{\infty}$ direct sum

$$
\mathbf{Y} \stackrel{\text { def }}{=} \mathbf{Y}^{m_{1}} \oplus_{\infty} \mathbf{Y}^{m_{2}} \oplus_{\infty} \ldots \oplus_{\infty} \mathbf{Y}^{m_{s+1}}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)
$$

Then $\|\cdot\|_{\ell_{\infty}^{n}} \leqslant\|\cdot\|_{\mathbf{Y}} \leqslant \max _{i \in\{1, \ldots, s+1\}} e^{C \log ^{*} m_{i}}\|\cdot\|_{\ell_{\infty}^{m_{i}}} \leqslant e^{C \log ^{*} n}\|\cdot\|_{\ell_{\infty}^{n}}$. We claim that $\frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \leqslant e^{O\left(\log ^{*} n\right)}$.
Since $B_{\mathbf{Y}}=B_{\mathbf{Y}^{m_{1}}} \times B_{\mathbf{Y}^{m_{2}}} \times \ldots \times B_{\mathbf{Y}^{m_{s+1}}}$, by an inductive application of Lemma 158 we have

$$
\frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{Y}}\right)} \leqslant\left(\sum_{i=1}^{s+1} \frac{\operatorname{MaxProj}\left(B_{\mathbf{Y}^{m_{i}}}\right)^{2}}{\operatorname{vol}_{m_{i}}\left(B_{\mathbf{Y}^{m_{i}}}\right)^{2}}\right)^{\frac{1}{2}} \leqslant\left(\sum_{i=1}^{s+1} e^{2 C \log ^{*} m_{i}}\right)^{\frac{1}{2}} \lesssim e^{C \log ^{*} n}
$$

where the first step uses Lemma 158 , the penultimate step is our assumption on $\mathbf{Y}^{m_{i}}$, and the final step has the following elementary justification. Recall that for every $i \in\{1, \ldots, s\}$ we have $m_{i+1} \leqslant\left(\log m_{i}\right)^{c}$, where $c>1$ is a universal constant. So, $m_{i+2} \leqslant c^{c}\left(\log \log m_{i}\right)^{c}$ for every $i \in\{1, \ldots, s-1\}$. Fix $n_{0} \in \mathbb{N}$ such that $c^{c}(\log \log n)^{c} \leqslant \log n$ for every $n \geqslant n_{0}$. Then, $m_{i+2} \leqslant \log m_{i}$ if $m_{i} \geqslant n_{0}$, hence $\log ^{*} m_{i+2} \leqslant \log ^{*} m_{i}-1$.

Let $i_{0}$ be the largest $i \in\{1, \ldots, s+1\}$ for which $m_{i}<n_{0}$. Then, $\log ^{*} m_{2 i} \leqslant \log ^{*} m_{2}-i \leqslant \log ^{*} n-i$ and $\log ^{*} m_{2 j+1} \leqslant \log ^{*} m_{1}-j \leqslant \log ^{*} n-j$ if $2 i, 2 j+1 \in\left\{1, \ldots, i_{0}-1\right\}$. Also, $\left|\left\{i_{0}, \ldots, s+1\right\}\right|=O(1)$. Consequently,

$$
\sum_{i=1}^{s+1} e^{2 C \log ^{*} m_{i}} \leqslant e^{2 C \log ^{*} n} \sum_{k=0}^{\infty} e^{-2 C k}+O(1) \lesssim e^{2 C \log ^{*} n} .
$$

Remark 164. A straightforward way to attempt to compute the surface area of the unit sphere of a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is to fix a direction $z \in S^{n-1}$ and consider $\partial B_{\mathbf{X}}$ as the union of the two graphs of the functions $\Psi_{z}^{\mathbf{X}}, \psi_{z}^{\mathbf{X}}: \operatorname{Proj}_{z^{\perp}}\left(B_{\mathbf{X}}\right) \rightarrow \mathbb{R}$ that are defined by setting $\Psi_{z}^{\mathbf{X}}(x)$ and $\psi_{z}^{\mathbf{X}}(x)$ for each $x \in \operatorname{Proj}_{z^{\perp}}\left(B_{\mathbf{X}}\right)$ to be, respectively, the largest and smallest $s \in \mathbb{R}$ for which $x+s z \in \partial B_{\mathbf{X}}$. We then have

$$
\begin{equation*}
\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right)=\int_{\operatorname{Proj}_{z^{\perp}}\left(B_{\mathbf{X}}\right)} \sqrt{1+\left\|\nabla \Psi_{Z}^{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}}^{2}} \mathrm{~d} x+\int_{\operatorname{Proj}_{z^{\prime}} \perp\left(B_{\mathbf{X}}\right)} \sqrt{1+\left\|\nabla \psi_{Z}^{\mathbf{X}}(x)\right\|_{\ell_{2}^{n}}^{2}} \mathrm{~d} x . \tag{427}
\end{equation*}
$$

When $\mathbf{X}=\ell_{p}^{n}$ for some $p \in(1, \infty)$ and $z=e_{n}$,

$$
\forall x \in \operatorname{Proj}_{e_{n}^{n}}\left(B_{\ell_{p}^{n}}\right)=B_{\ell_{p}^{n-1}}, \quad \Psi_{e_{n}}^{\ell_{p}^{n}}(x)=-\psi_{e_{n}}^{\ell_{p}^{n}}(x)=\left(1-\|x\|_{\ell_{p}^{n-1}}^{p}\right)^{\frac{1}{p}} .
$$

Therefore (427) becomes

$$
\frac{\operatorname{vol}_{n-1}\left(\partial B_{\ell_{p}^{n}}\right)}{\operatorname{vol}_{n-1}\left(B_{\ell_{p}^{n-1}}\right)}=2 f_{B_{\ell_{p}^{n-1}}}\left(1+\left(1-\|x\|_{\ell_{p}^{n-1}}^{p}\right)^{-\frac{2(p-1)}{p}} \sum_{i=1}^{n-1}\left|x_{i}\right|^{2(p-1)}\right)^{\frac{1}{2}} \mathrm{~d} x .
$$

By BGMN05, a point chosen from the normalized volume measure on $B_{\ell_{p}^{n-1}}$ is equidistributed with

$$
\left(\left|\mathrm{G}_{1}\right|^{p}+\ldots+\left|\mathrm{G}_{n-1}\right|^{p}+\mathrm{Z}\right)^{-\frac{1}{p}}\left(\mathrm{G}_{1}, \ldots, \mathrm{G}_{n-1}\right) \in \mathbb{R}^{n-1}
$$

where $\mathrm{G}_{1}, \ldots, \mathrm{G}_{n-1}, \mathrm{Z}$ are independent random variables, the density of $\mathrm{G}_{1}, \ldots, \mathrm{G}_{n-1}$ at $s \in \mathbb{R}$ is equal to $2 \Gamma(1+1 / p)^{-1} e^{-|s|^{p}}$ and the density of $Z$ at $t \in[0, \infty)$ is equal to $e^{-t}$. Consequently,

$$
\begin{equation*}
\frac{\operatorname{vol}_{n-1}\left(\partial B_{\ell_{p}^{n}}\right)}{\operatorname{vol}_{n-1}\left(B_{\ell_{p}^{n-1}}\right)}=2 \mathbb{E}\left[\left(1+\mathrm{Z}^{-\frac{2(p-1)}{p}} \sum_{i=1}^{n-1}\left|\mathrm{G}_{i}\right|^{2(p-1)}\right)^{\frac{1}{2}}\right] . \tag{428}
\end{equation*}
$$

Optimal estimates on moments such as the right hand side of (428) were derived (in greater generality) in (Nao07], using which one can quickly get asymptotically sharp bounds on the left hand side of (428). It is possible to implement this approach to get an alternative treatment of $\ell_{p}^{n}\left(\ell_{q}^{m}\right)$, though it is significantly more involved than the different way by which we proceeded above, and it becomes much more tedious and technically intricate when one aims to treat hierarchically nested $\ell_{p}$ norms as we did in Lemma 161. Nevertheless, an advantage of (427) is that it applies to normed spaces that do not have a product structure as in Lemma 156, which is helpful in other settings that we will study elsewhere.
6.2. Negatively correlated normed spaces. Our goal here is to further elucidate the role of symmetries in the context of the discussion in Section 1.6.2. Fix $n \in \mathbb{N}$ and $\gamma \geqslant 1$. Say that a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{x}}\right)$ is $\gamma$-negatively correlated if the standard scalar product $\langle\cdot, \cdot\rangle$ on $\mathbb{R}^{n}$ is invariant under its isometry group $\operatorname{Isom}(\mathbf{X})$, i.e., $\operatorname{Isom}(\mathbf{X}) \leqslant \mathrm{O}_{n}$, and there exists a Borel probability measure $\mu$ on Isom(X) such that

$$
\begin{equation*}
\forall x, y \in \mathbb{R}^{n}, \quad \int_{\operatorname{Isom}(\mathbf{X})}|\langle U x, y\rangle| \mathrm{d} \mu(U) \leqslant \frac{\gamma}{\sqrt{n}}\|x\|_{\ell_{2}^{n}}\|y\|_{\ell_{2}^{n}} . \tag{429}
\end{equation*}
$$

We were inspired to formulate this notion by the proof of Theorem 1.1 in [Sch89]. It is tailored for the purpose of bounding volumes of hyperplane projections of $B_{\mathbf{X}}$ from above in terms of the surface area of $\partial B_{\mathbf{X}}$, as exhibited by the following lemma which generalizes the reasoning in [Sch89.
Lemma 165. Fix $n \in \mathbb{N}$ and $\gamma \geqslant 1$. If $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is $\gamma$-negatively correlated, then

$$
\operatorname{MaxProj}\left(B_{\mathbf{X}}\right) \leqslant \frac{\gamma}{2 \sqrt{n}} \operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right)
$$

Proof. Recall that for every $y \in \partial B_{\mathbf{X}}$ at which $\partial B_{\mathbf{X}}$ is smooth we denote the unit outer normal to $\partial B_{\mathbf{X}}$ at $y$ by $N_{\mathbf{X}}(y) \in S^{n-1}$. By the Cauchy projection formula (30) for every $x \in S^{n-1}$ we have

$$
\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}}\left(B_{\mathbf{X}}\right)\right)=\frac{1}{2} \int_{\partial B_{\mathbf{X}}}\left|\left\langle x, N_{\mathbf{X}}(y)\right\rangle\right| \mathrm{d} y .
$$

Since every $U \in \operatorname{Isom}(\mathbf{X})$ is an orthogonal transformation and $N_{\mathbf{X}} \circ U^{*}=U^{*} \circ N_{\mathbf{X}}$ almost surely on $\partial B_{\mathbf{X}}$,

$$
\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}}\left(B_{\mathbf{X}}\right)\right)=\frac{1}{2} \int_{\partial B_{\mathbf{X}}}\left|\left\langle U x, N_{\mathbf{X}}(y)\right\rangle\right| \mathrm{d} y .
$$

By integrating this identity with respect to $\mu$, we therefore conclude that

$$
\operatorname{vol}_{n-1}\left(\operatorname{Proj}_{x^{\perp}}\left(B_{\mathbf{X}}\right)\right)=\frac{1}{2} \int_{\partial B_{\mathbf{X}}}\left(\int_{\operatorname{Isom}(\mathbf{X})}\left|\left\langle U x, N_{\mathbf{X}}(y)\right\rangle\right| \mathrm{d} \mu(U)\right) \mathrm{d} y \leqslant \frac{\gamma}{2 \sqrt{n}} \operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right)
$$

where we used 429) and the fact that $\|x\|_{\ell_{2}^{n}}=1$ and $\left\|N_{\mathbf{X}}(y)\right\|_{\ell_{2}^{n}}=1$ for almost every $y \in \partial B_{\mathbf{X}}$.
By substituting Lemma 165 into Theorem 75 and using 97, we get the following corollary.
Corollary 166. Fix $n \in \mathbb{N}$ and $\gamma \geqslant 1$. If $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is $\gamma$-negatively correlated, then

$$
\mathrm{e}(\mathbf{X}) \lesssim \operatorname{SEP}(\mathbf{X}) \leqslant 2 \gamma \frac{\operatorname{vol}_{n-1}\left(\partial B_{\mathbf{X}}\right) \operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right) \sqrt{n}}
$$

Corollary 166 generalizes Corollary 44 since any canonically positioned normed space is 1-negatively correlated. Indeed, suppose that $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ is canonically positioned. Recall that in Section 1.6.2 we denoted the Haar probability measure on $\operatorname{Isom}(\mathbf{X})$ by $h_{\mathbf{X}}$. Fix $x, y \in \mathbb{R}^{n}$. The distribution of the random vector $U x$ when $U$ is distributed according to $h_{\mathbf{X}}$ is $\operatorname{Isom}(\mathbf{X})$-invariant, and therefore it is isotropic. Hence,

$$
\begin{aligned}
\int_{\operatorname{Isom}(\mathbf{X})}|\langle U x, y\rangle| \mathrm{d} h_{\mathbf{X}}(U) & \leqslant\left(\int_{\operatorname{Isom}(\mathbf{X})}\langle U x, y\rangle^{2} \mathrm{~d} h_{\mathbf{X}}(U)\right)^{\frac{1}{2}} \\
& \stackrel{70}{=} \frac{\|y\|_{\ell_{2}^{n}}}{\sqrt{n}}\left(\int_{\operatorname{Isom}(\mathbf{X})}\|U x\|_{\ell_{2}^{n}}^{2} \mathrm{~d} h_{\mathbf{X}}(U)\right)^{\frac{1}{2}}=\frac{1}{\sqrt{n}}\|x\|_{\ell_{2}^{n}}\|y\|_{\ell_{2}^{n}},
\end{aligned}
$$

where the final step uses the fact that each $U \in \operatorname{Isom}(\mathbf{X})$ is an orthogonal transformation.
One way to achieve (429, which is close in spirit to the considerations in [Sch89], is when there are $\Gamma \subseteq\{-1,1\}^{n}$ and $\Pi \subseteq S_{n}$ such that $U_{\varepsilon, \pi} \in \operatorname{lsom}(\mathbf{X})$ for every $(\varepsilon, \pi) \in \Gamma \times \Pi$, where $U_{\varepsilon, \pi} \in G L_{n}(\mathbb{R})$ is given by

$$
\forall x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}, \quad U_{\varepsilon, \pi} x \stackrel{\text { def }}{=}\left(\varepsilon_{1} x_{\pi(1)}, \ldots \varepsilon_{n} x_{\pi(n)}\right)
$$

and also there are $\alpha, \beta>0$ such that

$$
\begin{equation*}
\forall w \in \mathbb{R}^{n}, \quad \frac{1}{|\Gamma|} \sum_{\varepsilon \in \Gamma}|\langle\varepsilon, w\rangle| \leqslant \alpha\|w\|_{\ell_{2}^{n}} \tag{430}
\end{equation*}
$$

and

$$
\begin{equation*}
\forall i, j \in\{1, \ldots, n\}, \quad|\{\pi \in \Pi: \pi(i)=j\}| \leqslant \beta \frac{|\Pi|}{n} \tag{431}
\end{equation*}
$$

Under these assumptions, $\mathbf{X}$ is $\gamma$-negatively correlated with $\gamma=\alpha \sqrt{\beta}$. Indeed, we can take $\mu$ in (429) to be the uniform distribution over the finite set $\left\{U_{\varepsilon, \pi}:(\varepsilon, \pi) \in \Gamma \times \Pi\right\} \subseteq \operatorname{Isom}(\mathbf{X})$, since every $x, y \in \mathbb{R}^{n}$ satisfy

$$
\begin{aligned}
& \frac{1}{|\Gamma \times \Pi|} \sum_{(\varepsilon, \pi) \in \Gamma \times \Pi}\left|\left\langle U_{\varepsilon, \pi} x, y\right\rangle\right| \stackrel{\boxed{430}}{\leqslant} \frac{1}{|\Pi|} \sum_{\pi \in \Pi} \alpha\left(\sum_{i=1}^{n}\left(x_{\pi(i)} y_{i}\right)^{2}\right)^{\frac{1}{2}} \leqslant \alpha\left(\sum_{i=1}^{n}\left(\frac{1}{|\Pi|} \sum_{\pi \in \Pi} x_{\pi(i)}^{2}\right) y_{i}^{2}\right)^{\frac{1}{2}} \\
&=\alpha\left(\sum_{i=1}^{n}\left(\frac{1}{|\Pi|} \sum_{j=1}^{n}|\{\pi \in \Pi: \pi(i)=j\}| x_{j}^{2}\right) y_{i}^{2}\right)^{\frac{1}{2}} \stackrel{431}{\leqslant} \frac{\alpha \sqrt{\beta}}{\sqrt{n}}\|x\|_{\ell_{2}^{n}}\|y\|_{\ell_{2}^{n}} .
\end{aligned}
$$

The condition (430) can be viewed as a negative correlation property of the coordinates of sign vectors that are chosen uniformly from $\Gamma$. The condition (431) roughly means that for each $i \in\{1, \ldots, n\}$ the sets $\{\pi \in \Pi: \pi(i)=1\}, \ldots,\{\pi \in \Pi: \pi(i)=n\}$ form an approximately equitable partition of $\Pi$. This holds with $\beta=1$ if $\Pi$ is a transitive subgroup of $S_{n}$. One could formulate weaker conditions that ensure the validity of the conclusion of Lemma 165 (e.g. considering bi-Lipschitz automorphisms of $\mathbf{X}$ rather than isometries of $\mathbf{X}$ ), and hence also the conclusion of Corollary 166, though we will not pursue this here as we expect that in concrete cases such issues should be easy to handle.
6.3. Volume ratio computations. Here we will present asymptotic evaluations of volume ratios of some normed spaces, for the purpose of plugging them into results that we stated in the Introdcution. Due to the large amount of knowledge on this topic that is available in the literature, we will only give a flavor of such applications. The main reference for the contents of this section is the valuable work [Sch82].

We will start by examining the iteratively nested $\ell_{p}$ products $\left\{\mathbf{X}_{k}\right\}_{k=0}^{\infty}$ of Lemma 152 , in the special case when the initial space $\mathbf{X}=\mathbf{X}_{0}$ is a canonically positioned normed space for which Conjecture 48 holds. Thus, we are fixing $\left\{n_{k}\right\}_{k=0}^{\infty} \subseteq \mathbb{N}$ and $\left\{p_{k}\right\}_{k=1}^{\infty} \subseteq[1, \infty]$, and assuming that $\mathbf{X}=\left(\mathbb{R}^{n_{0}},\|\cdot\|_{\mathbf{X}}\right)$ is a canonically positioned normed space satisfying Conjecture 48, i.e., 326) holds with $\alpha=O(1)$; the case $\mathbf{X}=\mathbb{R}$ is sufficiently rich for our present illustrative purposes, but one can also take $\mathbf{X}=\mathbf{E}$ to be any symmetric space, per Lemma 53. By Lemma 152 and Corollary 78 , if we define inductively

$$
\forall k \in \mathbb{N}, \quad \mathbf{X}_{k+1}=\ell_{p_{k}}^{n_{k}}\left(\mathbf{X}_{k}\right), \quad \text { where } \quad \mathbf{X}_{0}=\mathbf{X}
$$

then, because $\left\{\mathbf{X}_{k}\right\}_{k=1}^{\infty}$ are canonically positioned (they belong to the class of spaces in Remark 39),

$$
\begin{equation*}
\forall m \in \mathbb{N}, \quad \operatorname{SEP}\left(\mathbf{X}_{m}\right)=\operatorname{evr}\left(\mathbf{X}_{m}\right) \sqrt{\operatorname{dim}\left(\mathbf{X}_{m}\right)}=\operatorname{evr}\left(\mathbf{X}_{m}\right) \sqrt{n_{0} \cdots n_{m}} \tag{432}
\end{equation*}
$$

Let $\left\{\mathbf{H}_{k}\right\}_{k=0}^{\infty}$ be the sequence of Euclidean spaces that arise from the above construction with the same $\left\{n_{k}\right\}_{k=0}^{\infty} \subseteq \mathbb{N}$ but with $p_{k}=2$ for all $k \in \mathbb{N}$ and $\mathbf{X}=\ell_{2}^{n_{0}}$. Thus, for each $m \in \mathbb{N}$ the Euclidean space $\mathbf{H}_{m}$ can be identified naturally with $\ell_{2}^{n_{0} \cdots n_{m}}$. Under this identification, by a straightforward inductive application of Hölder's inequality and the fact that the $\ell_{p}$ norm deceases with $p$, the Löwner ellipsoid of $\mathbf{X}_{m}$ satisfies ${ }^{16}$

$$
\mathcal{L}_{\mathbf{X}_{m}} \subseteq\left(\prod_{k=1}^{m} n_{k}^{\max \left\{\frac{1}{2}-\frac{1}{p_{k}}, 0\right\}}\right)\left(\mathcal{L}_{\mathbf{X}}\right)^{n_{1} \cdots n_{m}} .
$$

Also, by Lemma 149 we have

$$
\operatorname{vol}_{n_{0} \cdots n_{m}}\left(B \mathbf{X}_{m}\right)^{\frac{1}{n_{0} \cdots n_{k}}}=\frac{\operatorname{vol}_{n_{0}}\left(B_{\mathbf{X}}\right)^{\frac{1}{n_{0}}}}{\prod_{k=1}^{m} n_{k}^{\frac{1}{p_{k}}}}
$$

These facts combine to give the following consequence of 432.

$$
\operatorname{SEP}\left(\mathbf{X}_{m}\right)=\operatorname{evr}(\mathbf{X}) \prod_{k=1}^{m} n_{k}^{\max \left\{\frac{1}{2}, \frac{1}{p_{k}}\right\}}
$$

In particular, when we take $X=\mathbb{R}$ and consider only two steps of the above iteration, we get the following asymptotic evaluation of the separation modulus of the $\ell_{p}^{n}\left(\ell_{q}^{m}\right)$ norm the space of $n$-by- $m$ matrices $M_{n \times m}(\mathbb{R})$ for any $n, m \in \mathbb{N}$ and $p, q \geqslant 1$; the case of square matrices was stated in the Introduction as (5).

$$
\begin{equation*}
\operatorname{SEP}\left(\ell_{p}^{n}\left(\ell_{q}^{m}\right)\right)=n^{\max \left\{\frac{1}{p}, \frac{1}{2}\right\}} m^{\max \left\{\frac{1}{q}, \frac{1}{2}\right\}}=\max \left\{\sqrt{n m}, m^{\frac{1}{q}} \sqrt{n}, n^{\frac{1}{p}} \sqrt{m}, n^{\frac{1}{p}} m^{\frac{1}{q}}\right\} \tag{433}
\end{equation*}
$$

Next, fix an integer $n \geqslant 2$ and let $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ be an unconditional normed space. Given $q \in[2, \infty]$ and $\Lambda \geqslant 1$, one says (see e.g. [LT79, Definitionl.f.4]) that $\mathbf{E}$ satisfies a lower $q$-estimate with constant $\Lambda$ if

[^11]for every $\left\{u_{k}\right\}_{k=1}^{\infty} \subseteq \mathbb{R}^{n}$ with pairwise disjoint supports we have
\[

$$
\begin{equation*}
\left(\sum_{k=1}^{\infty}\left\|u_{k}\right\|_{\mathbf{E}}^{q}\right)^{\frac{1}{q}} \leqslant \Lambda\left\|\sum_{k=1}^{\infty} u_{k}\right\|_{\mathbf{E}} \tag{434}
\end{equation*}
$$

\]

Note that by (323) this always holds with $\Lambda=1$ if $q=\infty$.
In concrete cases it is often mechanical to evaluate up to universal constant factors the minimum radius of a Euclidean ball that circumscribes $B_{\mathbf{X}}$, but it is always within a $O(\sqrt{\log n})$ factor of the expression

$$
\begin{equation*}
R_{\mathbf{E}} \stackrel{\text { def }}{=} \max _{\varnothing \neq S \subseteq\{1, \ldots, n\}}\left(\frac{\sqrt{|S|}}{\left\|\sum_{i \in S} e_{i}\right\|_{\mathbf{E}}}\right) . \tag{435}
\end{equation*}
$$

More precisely, if $\mathbf{E}$ satisfies a lower $q$-estimate with constant $\Lambda$, then

$$
\begin{equation*}
R_{\mathbf{E}} \leqslant \text { outradius }_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right) \lesssim \Lambda(\log n)^{\frac{1}{2}-\frac{1}{q}} R_{\mathbf{E}} \tag{436}
\end{equation*}
$$

The first inequality in (436) is immediate because $\left\|\sum_{i \in S} e_{i}\right\|_{\mathrm{E}}^{-1} \sum_{i \in S} e_{i} \in B_{\mathrm{E}}$ if $\varnothing \neq S \subseteq\{1, \ldots, n\}$. For a quick justification of the second inequality in (436), note that by homogeneity we may assume without loss of generality that $\left\|e_{i}\right\|_{\mathbf{E}} \geqslant 1$ for every $i \in \mathbb{N}$. Therefore, using (323) we see that if $x=\left(x_{1}, \ldots, x_{n}\right) \in B_{\mathbf{E}}$, then $\max _{i \in\{1, \ldots, n\}}\left|x_{i}\right| \leqslant 1$. Consequently, if we fix $x \in B_{\mathbf{E}}$ and denote for each $k \in \mathbb{N}$,

$$
\begin{equation*}
S_{k}=S_{k}(x) \stackrel{\text { def }}{=}\left\{i \in\{1, \ldots, n\}: \frac{1}{2^{k}}<\left|x_{i}\right| \leqslant \frac{1}{2^{k-1}}\right\} \tag{437}
\end{equation*}
$$

then the sets $\left\{S_{k}\right\}_{k=1}^{\infty}$ are a partition of $\{1, \ldots, n\}$ and in particular $\sum_{k=1}^{\infty}\left|S_{k}\right|=n$. Next,

$$
\begin{equation*}
\Lambda R_{\mathbf{E}} \geqslant \Lambda R_{\mathbf{E}}\|x\|_{\mathbf{E}} \geqslant R_{\mathbf{E}}\left(\sum_{k=1}^{\infty}\left\|\sum_{i \in S_{k}} x_{i} e_{i}\right\|_{\mathbf{E}}^{q}\right)^{\frac{1}{q}} \geqslant\left(\sum_{k=1}^{\infty} R_{\mathbf{E}}^{q}\left\|\sum_{i \in S_{k}} \frac{1}{2^{k}} e_{i}\right\|_{\mathbf{E}}^{q}\right)^{\frac{1}{q}} \geqslant\left(\sum_{k=1}^{\infty} \frac{\left|S_{k}\right|^{\frac{q}{2}}}{2^{k q}}\right)^{\frac{1}{q}} \tag{438}
\end{equation*}
$$

where the second step of (438) uses (434), the penultimate step of 438) uses (323) and 437), and the final step of 438) uses 435. Now, for every $0<\theta<1$ we have

$$
\begin{align*}
\|x\|_{\ell_{2}^{n}} & =\left(\sum_{k=1}^{\infty} \sum_{i \in S_{k}} x_{i}^{2}\right)^{\frac{1}{2}} \leqslant\left(\sum_{k=1}^{\infty} \frac{\left|S_{k}\right|}{2^{2(k-1)}}\right)^{\frac{1}{2}}=2\left(\sum_{k=1}^{\infty} \frac{\left|S_{k}\right|^{1-\theta}}{2^{2 k(1-\theta)}}\left|S_{k}\right|^{\theta} 2^{-2 k \theta}\right)^{\frac{1}{2}} \\
& \leqslant 2\left(\sum_{k=1}^{\infty} \frac{\left|S_{k}\right|^{\frac{q}{2}}}{2^{2 k q}}\right)^{\frac{1-\theta}{q}}\left(\sum_{k=1}^{\infty}\left|S_{k}\right|\right)^{\frac{\theta}{2}}\left(\sum_{k=1}^{\infty} 2^{-\frac{2 k q \theta}{(q-2)(1-\theta)}}\right)^{\left(\frac{1}{2}-\frac{1}{q}\right)(1-\theta)} \lesssim\left(\Lambda R_{\mathbf{E}}\right)^{1-\theta} n^{\frac{\theta}{2}} \theta^{-\left(\frac{1}{2}-\frac{1}{q}\right)}, \tag{439}
\end{align*}
$$

where the second step of (439) uses (437), the penultimate step of (439) uses the trilinear Hölder inequality with exponents $1 / \theta, q /(2(1-\theta))$ and $1 /((1-2 / q)(1-\theta))$, and the final step of 439 ) uses 438), the fact that $\sum_{k=1}^{\infty}\left|S_{k}\right|=n$, and elementary calculus. By choosing $\theta=1 / \log n$ in (439), we get 436).

By the Lozanovskiĭ factorization theorem Loz69] there exist $w_{1}, \ldots, w_{n}>0$ such that

$$
\begin{equation*}
\left\|\sum_{i=1}^{n} w_{i} e_{i}\right\|_{\mathbf{E}}=\left\|\sum_{i=1}^{n} \frac{1}{w_{i}} e_{i}\right\|_{\mathbf{E}^{*}}=\sqrt{n} \tag{440}
\end{equation*}
$$

We will call any $w_{1}, \ldots, w_{n}>0$ that satisfy (440) Lozanovskiĭ weights for $\mathbf{E}$. They can be found by maximizing the concave function $w \mapsto \sum_{i=1}^{n} \log w_{i}$ over $w \in B_{\mathbf{E}}$ (see also e.g. [Pis89, Chapter 3]), which can be done efficiently if $\mathbf{E}$ is given by an efficient oracle; their existence can also be established nonconstructively using the Brouwer fixed point theorem [JR76]. By [Sch82, Lemma 1.2] (note that we are using a different normalization of the weights than in Sch82]),

$$
\begin{equation*}
\operatorname{vol}_{n}\left(B_{\mathbf{E}}\right)^{\frac{1}{n}}=\frac{\left(w_{1} \cdots w_{n}\right)^{\frac{1}{n}}}{\sqrt{n}} \tag{441}
\end{equation*}
$$

By combining (436) and (441, we get the following lemma.

Lemma 167. Fix an integer $n \geqslant 2$ and let $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ be an unconditional normed space. Suppose that $\mathbf{E}$ satisfies a lower $q$-estimate with constant $\Lambda$ for some $q \geqslant 2$ and $\Lambda \geqslant 1$. Then,

$$
\operatorname{evr}(\mathbf{E}) \lesssim \frac{\max _{\varnothing \neq S \subseteq\{1, \ldots, n\}}\left(\frac{\sqrt{|S|}}{\left\|\sum_{i \in S} e_{i}\right\|_{\mathbf{E}}}\right)}{\sqrt[n]{w_{1} \cdots w_{n}}} \Lambda(\log n)^{\frac{1}{2}-\frac{1}{q}},
$$

for any Lozanovskiĭ weights $w_{1}, \ldots, w_{n}>0$ for $\mathbf{E}$. If the Löwner ellipsoid of $\mathbf{E}$ is a multiple of $B_{\ell_{2}^{n}}$, then

$$
\frac{\max _{\varnothing \nexists S \subseteq\{1, \ldots, n\}}\left(\frac{\sqrt{|S|}}{\left\|\sum_{i \in S} e_{i}\right\|_{\mathbf{E}}}\right)}{\sqrt[n]{w_{1} \cdots w_{n}}} \lesssim \operatorname{evr}(\mathbf{E}) \lesssim \frac{\max _{\varnothing \neq S \subseteq\{1, \ldots, n\}}\left(\frac{\sqrt{|S|}}{\left\|\sum_{i \in S} e_{i}\right\|_{\mathbf{E}}}\right)}{\sqrt[n]{w_{1} \cdots w_{n}}} \Lambda(\log n)^{\frac{1}{2}-\frac{1}{q}} .
$$

The following corollary is a consequence of Lemma 167 because if $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ is a normed space that satisfies the assumptions of Lemma 52 (in particular, $\mathbf{E}$ is unconditional), then by Lemma 151

$$
w_{1}=w_{2}=\ldots=w_{n}=\frac{\sqrt{n}}{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}}
$$

are Lozanovskiĭ weights for $\mathbf{E}$.
Corollary 168. If $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ a normed space that satisfies the assumptions of Lemma 52 , then

$$
\frac{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}}{\sqrt{n}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) \lesssim \operatorname{evr}(\mathbf{E}) \lesssim \frac{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}}{\sqrt{n}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) \sqrt{\log n}
$$

Hence, by Corollary 78 we have

$$
\left\|e_{1}+\ldots+e_{n}\right\|_{\mathrm{E}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) \lesssim \operatorname{SEP}(\mathbf{E}) \lesssim\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) \sqrt{\log n}
$$

More succinctly, this can be written in the following form, which we already stated in Corollary 6 .

$$
\begin{equation*}
\operatorname{SEP}(\mathbf{E})=\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) n^{o(1)} \tag{442}
\end{equation*}
$$

By Sch82, Proposition 2.2], the unitary ideal of any symmetric normed space $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ satisfies

$$
\begin{equation*}
\operatorname{vr}\left(S_{E}\right)=\operatorname{vr}(\mathbf{E}) . \tag{443}
\end{equation*}
$$

This implies that

$$
\begin{equation*}
\operatorname{evr}\left(S_{\mathbf{E}}\right)=\operatorname{evr}(\mathbf{E}), \tag{444}
\end{equation*}
$$

by (72) combined with $\mathrm{S}_{\mathbf{E}}^{*}=\mathrm{S}_{\mathbf{E}^{*}}$, though a straightforward adjustment of the proof of (443) in [Sch82] yields (444) directly, without using the much deeper result (72). We therefore have the following corollary.

Corollary 169. If $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$ is a symmetric normed space, then

$$
\frac{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}}{\sqrt{n}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) \lesssim \operatorname{evr}\left(\mathrm{S}_{\mathbf{E}}\right) \lesssim \frac{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}}{\sqrt{n}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) \sqrt{\log n .}
$$

Hence, by Corollary 78 we have
$\left\|e_{1}+\ldots+e_{n}\right\|_{\mathrm{E}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) \sqrt{n} \lesssim \operatorname{SEP}\left(\mathrm{~S}_{\mathrm{E}}\right) \lesssim\left\|e_{1}+\ldots+e_{n}\right\|_{\mathrm{E}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) \sqrt{n} \log n$,
More succinctly, this can be written in the following form, which we already stated in Corollary 6 ,

$$
\begin{equation*}
\operatorname{SEP}\left(\mathrm{S}_{\mathbf{E}}\right)=\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}\left(\max _{k \in\{1, \ldots, n\}} \frac{\sqrt{k}}{\left\|e_{1}+\ldots+e_{k}\right\|_{\mathbf{E}}}\right) n^{\frac{1}{2}+o(1)} \tag{445}
\end{equation*}
$$

Remark 170. In the above discussion, as well as in the ensuing treatment of tensor products, we prefer to consider square matrices rather than rectangular matrices because the setting of square matrices exhibits all of the key issues while being notationally simpler. Nevertheless, there are two places in which we do need to work with rectangular matrices, namely the above proof of Proposition 163 and the proof of the first inequality in (118). For the latter, fix $p \geqslant 1$ and $n, m \in \mathbb{N}$. As in the proof of Theorem 76, denote the Schatten-von Neumann trace class on the $n$-by- $m$ real matrices $\mathrm{M}_{n \times m}(\mathbb{R})$ by $\mathrm{S}_{p}^{n \times m}$; recall $(120)$. The following asymptotic identity implies (121) (recall that in the setting of (121) we have $r \in\{1, \ldots, n\}$ ).

$$
\begin{equation*}
\operatorname{evr}\left(\mathrm{S}_{p}^{n \times m}\right)=(\min \{n, m\})^{\max \left\{\frac{1}{p}-\frac{1}{2}, 0\right\}} \tag{446}
\end{equation*}
$$

Volumes of unit balls of Schatten-von Neumann trace classes have been satisfactorily estimated in the literature, starting with STJ80] and the comprehensive work [Sch82], through the more precise asymptotics in SR84, KPT20]. Unfortunately, all of these works dealt only with square matrices. Nevertheless, these references could be mechanically adjusted to treat rectangular matrices as well. Since (446) does not seem to have been stated in the literature, we will next sketch its derivation by mimicking the reasoning of [Sch82], though the more precise statements of [SR84, KPT20] could be derived as well via similarly straighforward modifications of the known proofs for square matrices. We claim that

$$
\begin{equation*}
\operatorname{vol}_{n m}\left(B_{\mathrm{S}_{p}^{n \times m}}\right)^{\frac{1}{n m}}=\frac{1}{(\min \{n, m\})^{\frac{1}{p}} \sqrt{\max \{n, m\}}} \tag{447}
\end{equation*}
$$

(447) gives (446) since $S_{p}^{n \times m}$ is canonically positioned, so by Hölder's inequality its Löwner ellipsoid is

$$
\mathcal{L}_{\mathrm{S}_{p}^{n \times m}}=(\min \{n, m\})^{\max \left\{\frac{1}{2}-\frac{1}{p}, 0\right\}_{\mathrm{S}_{2}^{n \times m}}} .
$$

To prove 447, note first that it follows from its special case $p=\infty$. Indeed, as $\mathrm{S}_{1}^{n \times m}=\left(\mathrm{S}_{\infty}^{n \times m}\right)^{*}$, by the Blaschke-Santaló inequality [Bla17, San49] and the Bourgain-Milman inequality [BM87] the case $p=1$ of (447) follows from its case $p=\infty$. Now, 447) follows in full generality since by Hölder's inequality.

$$
\frac{1}{(\min \{n, m\})^{\frac{1}{p}}} B_{\mathrm{S}_{\infty}^{n \times m}} \subseteq B_{\mathrm{S}_{p}^{n \times m}} \subseteq(\min \{n, m\})^{1-\frac{1}{p}} B_{\mathrm{S}_{1}^{n \times m}}
$$

The upper bound $\operatorname{vol}_{n m}\left(B_{\mathrm{S}_{\infty}^{n \times m}}\right)^{1 /(n m)} \lesssim 1 / \sqrt{\max \{n, m\}}$ follows from $B_{\mathrm{S}_{\infty}^{n \times m}} \subseteq \sqrt{\min \{n, m\}} B_{\mathrm{S}_{2}^{n \times m}}$. For the matching lower bound, if $\left\{\varepsilon_{i j}\right\}_{i, j \in \mathbb{N}}$ are i.i.d. Bernoulli random variables, then by [BGN75, Theorem 1],

$$
\mathbb{E}\left[\left\|\sum_{i=1}^{n} \sum_{j=1}^{m} \varepsilon_{i j} e_{i} \otimes e_{j}\right\|_{\mathrm{S}_{\infty}^{n \times m}}\right] \lesssim \sqrt{\max \{n, m\}}
$$

This implies the lower bound $\operatorname{vol}_{n m}\left(B_{S_{\infty}^{n \times m}}\right)^{1 /(n m)} \gtrsim 1 / \sqrt{\max \{n, m\}}$ by [Sch82, Lemma 1.5].
Proof of Lemma53. By equation (2.2) in Sch82 we have

$$
\begin{equation*}
\operatorname{vol}_{n^{2}}\left(B_{\mathrm{S}_{\mathrm{E}}}\right)^{\frac{1}{n^{2}}}=\frac{1}{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathrm{E}} \sqrt{n}} \tag{448}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
\forall q \geqslant 1, \quad \operatorname{vol}_{n^{2}}\left(B_{\mathrm{S}_{q}^{n}}\right)^{\frac{1}{n^{2}}}=\frac{1}{n^{\frac{1}{2}+\frac{1}{q}}} \tag{449}
\end{equation*}
$$

Because $S_{q}^{n}$ is canonically positioned (it belongs to the class of spaces in Example 39), and hence it is in its minimum surface area position, by combining GP99, Proposition 3.1] and (54) we see that

$$
\begin{equation*}
\frac{\operatorname{vol}_{n^{2}-1}\left(\partial B_{\mathrm{S}_{q}^{n}}\right)}{\operatorname{vol}_{n^{2}}\left(B_{\mathrm{S}_{q}^{n}}\right)}=\frac{n \operatorname{MaxProj}\left(B_{\mathrm{S}_{q}^{n}}\right)}{\operatorname{vol}_{n^{2}}\left(B_{\mathrm{S}_{q}^{n}}\right)} \stackrel{43}{=} n^{\frac{3}{2}+\frac{1}{q}} \sqrt{\min \{q, n\}} \tag{450}
\end{equation*}
$$

Consequently,

$$
\begin{equation*}
\operatorname{iq}\left(B_{S_{q}^{n}}\right)=n \frac{\operatorname{vol}_{n^{2}-1}\left(\partial B_{q}^{n}\right)}{\operatorname{vol}_{n^{2}}\left(B_{S_{q}^{n}}\right)} \operatorname{vol}_{n^{2}}\left(B_{S_{q}^{n}}\right)^{\frac{1}{n^{2}}} \xlongequal\left[(449]{\wedge}=\frac{\sqrt[450]{ }}{n^{\frac{3}{2}+\frac{1}{q}} \sqrt{\min \{q, n\}}} n^{\frac{1}{2}+\frac{1}{q}}\right) n \sqrt{\min \{q, n\}} . \tag{451}
\end{equation*}
$$

Because by (323) we have $\|x\|_{\mathbf{E}} \leqslant\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}\|x\|_{\ell_{\infty}^{n}}$ for every $x \in \mathbb{R}^{n}$, every matrix $A \in \mathrm{M}_{n}(\mathbb{R})$ satisfies $\|A\|_{\mathrm{S}_{\mathbf{E}}} \leqslant\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}\|A\|_{\mathrm{S}_{\infty}^{n}} \leqslant\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}\|A\|_{\mathrm{S}_{q}^{n}}$. Consequently,

$$
\begin{equation*}
\frac{1}{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}} B_{\mathrm{S}_{q}^{n}} \subseteq B_{\mathrm{S}_{\mathbf{E}}} \tag{452}
\end{equation*}
$$

Moreover,

$$
\operatorname{iq}\left(\frac{1}{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}} B_{\mathrm{S}_{q}^{n}}\right)=\mathrm{iq}\left(B_{\mathrm{S}_{q}^{n}}\right) \stackrel{451}{=} n \sqrt{\min \{q, n\}}
$$

and

$$
\operatorname{vol}_{n^{2}}\left(\frac{1}{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}} B_{\mathrm{S}_{q}^{n}}\right)^{\frac{1}{n^{2}}} \stackrel{\boxed{449}}{=} \frac{1}{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}} n^{\frac{1}{2}+\frac{1}{q}}} \stackrel{448}{=} \frac{\operatorname{vol}_{n^{2}}\left(B_{\mathrm{S}_{\mathbf{E}}}\right)^{\frac{1}{n^{2}}}}{n^{\frac{1}{q}}}
$$

By choosing $q=\log n$ we get (81) for the normed space $\mathbf{Y}$ whose unit ball is the left hand side of 452).
Remark 171. An inspection of the proof of Lemma 53 reveals that if Conjecture 48 holds for $S_{\infty}^{n}$, then also Conjecture 48 holds for $S_{\mathbf{E}}$ for any symmetric normed space $\mathbf{E}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{E}}\right)$. Indeed, we would then take $\mathbf{Y}^{\prime}=\left(\mathrm{M}_{n}(\mathbb{R}),\|\cdot\|_{\mathbf{Y}^{\prime}}\right)$ to be the normed space whose unit ball is

$$
B_{\mathbf{Y}^{\prime}}=\frac{1}{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}} \operatorname{Ch} S_{\infty}^{n}=\frac{1}{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}} \mathrm{S}_{\chi \ell_{\infty}^{n}}
$$

where we recall Corollary 42. If Conjecture 48 holds for $S_{\infty}^{n}$, then $n=\mathrm{iq}\left(\operatorname{Ch} S_{\infty}^{n}\right)=\mathrm{iq}\left(B_{\mathbf{Y}^{\prime}}\right)$, and also

$$
\operatorname{vol}_{n^{2}}\left(\operatorname{Ch} S_{\infty}^{n}\right)^{\frac{1}{n^{2}}}=\operatorname{vol}_{n^{2}}\left(S_{\infty}^{n}\right)^{\frac{1}{n^{2}}} \stackrel{449}{=} \frac{1}{\sqrt{n}}
$$

from which we see that

$$
\operatorname{vol}_{n^{2}}\left(B_{\mathbf{Y}^{\prime}}\right)^{\frac{1}{n^{2}}}=\frac{\operatorname{vol}_{n^{2}}\left(\operatorname{Ch} S_{\infty}^{n}\right)^{\frac{1}{n^{2}}}}{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}}=\frac{1}{\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}} \sqrt{n}} \stackrel{448}{=} \operatorname{vol}_{n^{2}}\left(B_{\mathrm{S}_{\mathbf{E}}}\right)^{\frac{1}{n^{2}}}
$$

This proves Conjecture 48 for $\mathrm{S}_{\mathbf{E}}$. Note in passing that this also implies that

$$
\frac{1}{\sqrt{n}}=\operatorname{vol}_{n^{2}}\left(S_{\chi \ell_{\infty}^{n}}\right)^{\frac{1}{n^{2}}} \stackrel{448}{=} \frac{1}{\left\|e_{1}+\ldots+e_{n}\right\|_{\chi \ell_{\infty}^{n}} \sqrt{n}}
$$

Therefore, if Conjecture 48 holds for $S_{\infty}^{n}$, then $\left\|e_{1}+\ldots+e_{n}\right\|_{\chi \ell_{\infty}^{n}}=1$. More generally, by mimicking the above reasoning we deduce that if Conjecture 48 holds for $\mathrm{S}_{\mathbf{E}}$, then $\left\|e_{1}+\ldots+e_{n}\right\|_{\chi \mathbf{E}}=\left\|e_{1}+\ldots+e_{n}\right\|_{\mathbf{E}}$, which would be a modest step towards Problem 43 .

Fix $n \in \mathbb{N}$ and $p, q \geqslant 1$. We claim that the volume ratio of the projective tensor product $\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}$ satisfies

$$
\begin{equation*}
\operatorname{vr}\left(\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}\right)=\Phi_{p, q}(n) \tag{453}
\end{equation*}
$$

where

$$
\Phi_{p, q}(n) \stackrel{\text { def }}{=} \begin{cases}1 & \text { if } 1 \leqslant p, q \leqslant 2  \tag{454}\\ n^{\frac{1}{2}-\frac{1}{p}} & \text { if } q \leqslant 2 \leqslant p \leqslant \frac{q}{q-1} \\ n^{\frac{1}{q}-\frac{1}{2}} & \text { if } q \leqslant 2 \leqslant \frac{q-1}{q-1} \leqslant p \\ n^{\frac{1}{2}-\frac{1}{q}} & \text { if } p \leqslant 2 \leqslant q \leqslant \frac{p}{p-1} \\ n^{\frac{1}{p}-\frac{1}{2}} & \text { if } p \leqslant 2 \leqslant \frac{p}{p-1} \leqslant q \\ 1 & \text { if } p, q \geqslant 2 \text { and } \frac{1}{p}+\frac{1}{q} \geqslant \frac{1}{2} \\ n^{\frac{1}{2}-\frac{1}{p}-\frac{1}{q}} & \text { if } \frac{1}{p}+\frac{1}{q} \leqslant \frac{1}{2}\end{cases}
$$

Assuming (454) for the moment, by substituting it into Theorem 3 we get that

$$
\operatorname{SEP}\left(\ell_{p}^{n} \check{\otimes} \ell_{q}^{n}\right) \gtrsim n \operatorname{vr}\left(\left(\ell_{p}^{n} \check{\otimes} \ell_{q}^{n}\right)^{*}\right)=n \operatorname{vr}\left(\ell_{p^{*}}^{n} \hat{\otimes} \ell_{q^{*}}^{n}\right) \approx n \Phi_{p^{*}, q^{*}}(n)= \begin{cases}n & \text { if } p, q \geqslant 2, \\ n^{\frac{1}{p}+\frac{1}{2}} & \text { if } \frac{q}{q-1} \leqslant p \leqslant 2 \leqslant q \\ n^{\frac{3}{2}-\frac{1}{q}} & \text { if } p \leqslant \frac{q}{q-1} \leqslant 2 \leqslant q \\ n^{\frac{1}{q}+\frac{1}{2}} & \text { if } \frac{p}{p-1} \leqslant q \leqslant 2 \leqslant p \\ n^{\frac{3}{2}-\frac{1}{p}} & \text { if } q \leqslant \frac{p}{p-1} \leqslant 2 \leqslant p \\ n & \text { if } p, q \leqslant 2 \text { and } \frac{1}{p}+\frac{1}{q} \leqslant \frac{3}{2} \\ n^{\frac{1}{p}+\frac{1}{q}-\frac{1}{2}} & \text { if } \frac{1}{p}+\frac{1}{q} \geqslant \frac{3}{2}\end{cases}
$$

Since for any two normed spaces $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ and $\mathbf{Y}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Y}}\right)$ the space of operators from $\mathbf{X}^{*}$ to $\mathbf{Y}$ is isometric to the injective tensor product $\mathbf{X}^{*} \ddot{\otimes} \mathbf{Y}$ (see e.g. [DFS08]), we get from this that

$$
\operatorname{SEP}\left(\mathrm{M}_{n}(\mathbb{R}),\|\cdot\|_{\ell_{p}^{n} \rightarrow \ell_{q}^{n}}\right)=\operatorname{SEP}\left(\ell_{p^{*}}^{n} \check{\otimes} \ell_{q}^{n}\right) \gtrsim \begin{cases}n & \text { if } p \leqslant 2 \leqslant q  \tag{455}\\ n^{\frac{3}{2}-\frac{1}{p}} & \text { if } 2 \leqslant p \leqslant q \\ n^{\frac{3}{2}-\frac{1}{q}} & \text { if } 2 \leqslant q \leqslant p \\ n^{\frac{1}{q}+\frac{1}{2}} & \text { if } p \leqslant q \leqslant 2 \\ n^{\frac{1}{p}+\frac{1}{2}} & \text { if } q \leqslant p \leqslant 2 \\ n & \text { if } \frac{2 p}{p+2} \leqslant q \leqslant 2 \leqslant p \\ n^{\frac{1}{q}-\frac{1}{p}+\frac{1}{2}} & \text { if } q \leqslant \frac{2 p}{p+2}\end{cases}
$$

Observe that the rightmost quantity in (455) coincides with the right hand side of (14). Since $\ell_{p}^{n} \check{\otimes} \ell_{q}^{n}$ belongs to the class of spaces in Remark 39, a positive answer to Conjecture 11 for $\ell_{p}^{n} \check{\otimes} \ell_{q}^{n}$ would imply the following asymptotic evaluation of $\operatorname{SEP}\left(\ell_{p}^{n} \check{\otimes} \ell_{q}^{n}\right)$, which is equivalent to (14).

$$
\operatorname{SEP}\left(\ell_{p}^{n} \check{\otimes} \ell_{q}^{n}\right)= \begin{cases}n & \text { if } p, q \geqslant 2, \\ n^{\frac{1}{2}+\frac{1}{p}} & \text { if } \frac{q}{q-1} \leqslant p \leqslant 2 \leqslant q \\ n^{\frac{3}{2}-\frac{1}{q}} & \text { if } p \leqslant \frac{q}{q-1} \leqslant 2 \leqslant q, \\ n^{\frac{1}{2}+\frac{1}{q}} & \text { if } \frac{p}{p-1} \leqslant q \leqslant 2 \leqslant p \\ n^{\frac{3}{2}-\frac{1}{p}} & \text { if } q \leqslant \frac{p}{p-1} \leqslant 2 \leqslant p \\ n & \text { if } p, q \leqslant 2 \text { and } \frac{1}{p}+\frac{1}{q} \leqslant \frac{3}{2} \\ n^{\frac{1}{p}+\frac{1}{q}-\frac{1}{2}} & \text { if } \frac{1}{p}+\frac{1}{q} \geqslant \frac{3}{2}\end{cases}
$$

Furthermore, by Theorem79 the leftmost quantity in (455) is bounded from above by $O(\log n)$ times the rightmost quantity in (455), thus implying the fourth bullet point of Corollary 4 ,

The asymptotic evaluation (453) of $\operatorname{vr}\left(\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}\right)$ was proved in Sch82] up to constant factors that depend on $p, q$, namely [Sch82, Theorem 3.1] states that

$$
\begin{equation*}
\forall p, q>1, \quad \operatorname{vr}\left(\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}\right)=p, q \Phi_{p, q}(n) \tag{456}
\end{equation*}
$$

If $2 \in\{p, q\}$ and also $\min \{p, q\} \leqslant 2$, then (456) is due to Szarek and Tomczak-Jaegermann [STJ80]. More recently, Defant and Michels [DM06] generalized (456) to projective tensor products of symmetric normed spaces that are either 2-convex or 2-concave. The proof of 456] in [Sch82] yields constants that degenerate as $\min \{p, q\}$ tends to 1 . We will therefore next improve the reasoning in [Sch82] to get 453].

Lemma 172. Fix $n \in \mathbb{N}$ and $p, q \geqslant 1$. Let $\left\{\varepsilon_{i j}\right\}, j \in\{1, \ldots, n\}$ be i.i.d. Bernoulli random variables (namely, they are independent and each of them is uniformly distributed over $\{-1,1\})$. Then,

$$
\mathbb{E}\left[\left\|\sum_{i=1}^{n} \sum_{j=1}^{n} \varepsilon_{i j} e_{i} \otimes e_{j}\right\|_{\ell_{p}^{n} \check{\otimes} \ell_{q}^{n}}\right]=n^{\beta(p, q)} \stackrel{\operatorname{def}}{=}\left\{\begin{array}{lll}
n^{\frac{1}{p}+\frac{1}{q}-\frac{1}{2}} & \text { if } & \max \{p, q\} \leqslant 2,  \tag{457}\\
n^{\frac{1}{\min (p, q)}} & \text { if } & \max \{p, q\} \geqslant 2 .
\end{array}\right.
$$

Citing the work [Che78] of Chevet, a version of Lemma 172 appears as Lemma 2.3 in [Sch82], except that in [Sch82, Lemma 2.3] the implicit constants in (457) depend on $p, q$. An inspection of the proof of (456) in [Sch82] reveals that this is the only source of the dependence of the constants on $p, q$ (in fact, for this purpose [Sch82] only needs half of (457], namely to bound from above its left hand side by its right hand side). Specifically, all of the steps within [Sch82] incur only a loss of a universal constant factor, and the proof of (456] in [Sch82] also appeals to inequalities in the earlier work [Sch78] of Schütt, as well a classical inequality of Hardy and Littlewood [HL34]; all of the constants in these cited inequalities are universal. Therefore, (453) will be established after we prove Lemma 172 .
Proof of Lemma 172. Denote the random matrix whose $(i, j)$ entry is $\varepsilon_{i j}$ by $\mathcal{\varepsilon} \in \mathrm{M}_{n}(\mathbb{R})$. Then, the goal is

$$
\begin{equation*}
\mathbb{E}\left[\|\mathcal{E}\|_{\rho_{p^{*}}^{n} \rightarrow \ell_{q}^{n}}\right]=n^{\beta(p, q)} . \tag{458}
\end{equation*}
$$

In fact, the lower bound on the expected norm in (458) holds always, i.e., for a universal constant $c>0$,

$$
\begin{equation*}
\forall A \in \mathrm{M}_{n}(\{-1,1\}), \quad\|A\|_{\rho_{p^{*}}^{n} \rightarrow \ell_{q}^{n}} \geqslant c n^{\beta(p, q)} . \tag{459}
\end{equation*}
$$

A justification of (459) appears in the proof of Proposition 3.2 of Bennett's work Ben77] (specifically, see the reasoning immediately after inequality (15) in (Ben77), where it is explained that we can take $c=1$ if $\min \left\{p^{*}, q\right\} \geqslant 2$ or $\max \left\{p^{*}, q\right\} \leqslant 2$, and that we can take $c=1 / \sqrt{2}$ otherwise.

Next, let $\left\{\mathrm{g}_{i j}\right\}_{i, j \in\{1, \ldots, n\}}$ be i.i.d. standard Gaussian random variables. By [Che78, Lemme 3.1],

$$
\begin{equation*}
\mathbb{E}\left[\left\|\sum_{i=1}^{n} \sum_{j=1}^{n} \mathrm{~g}_{i j} e_{i} \otimes e_{j}\right\|_{\ell_{p}^{n} \stackrel{\otimes}{\otimes}{ }_{q}^{n}}\right]=n^{\max \left\{\frac{1}{p}+\frac{1}{q}-\frac{1}{2}, \frac{1}{p}\right\}} \sqrt{p}+n^{\max \left\{\frac{1}{p}+\frac{1}{q}-\frac{1}{2}, \frac{1}{q}\right\}} \sqrt{q} . \tag{460}
\end{equation*}
$$

Consequently,

$$
\begin{equation*}
\mathbb{E}\left[\left\|\sum_{i=1}^{n} \sum_{j=1}^{n} \varepsilon_{i j} e_{i} \otimes e_{j}\right\|_{\ell_{p}^{n} \check{\otimes} \ell_{q}^{n}}\right] \leqslant \sqrt{\frac{\pi}{2}} \mathbb{E}\left[\left\|\sum_{i=1}^{n} \sum_{j=1}^{n} g_{i j} e_{i} \otimes e_{j}\right\|_{\ell_{p}^{n} \check{\otimes} \ell_{q}^{n}}\right] \lesssim n^{\beta(p, q)} \sqrt{\max \{p, q\}}, \tag{461}
\end{equation*}
$$

where the first step of (461) is a standard comparison between Rademacher and Gaussian averages (a quick consequence of Jensen's inequality; e.g. (MP76]) and final step of (461) uses (460). This proves the desired bound (457) when $\max \{p, q\} \leqslant 2$, so suppose from now on that $\max \{p, q\} \geqslant 2$.

It suffices to treat the case $p \geqslant 2$. Indeed, if $p \leqslant 2$, then $q \geqslant 2$ since $\max \{p, q\} \geqslant 2$, so by the duality

$$
\|\mathcal{E}\|_{\ell_{p^{*}}^{n} \rightarrow \ell_{q}^{n}}=\left\|\mathcal{E}^{*}\right\|_{q_{q^{*}}^{n} \rightarrow \ell_{p}^{n}}
$$

and the fact that the transpose $\mathcal{E}^{*}$ has the same distribution as $\mathcal{E}$, the case $p \leqslant 2$ follows from the case $p \geqslant 2$. It also suffices to treat the case $q \leqslant p$ because if $q \geqslant p$, then $\|\cdot\|_{\ell_{q}^{n}} \leqslant\|\cdot\|_{\ell_{p}^{n}}$ point-wise, and therefore

$$
\|\mathcal{E}\|_{\ell_{p^{*}}^{n}} \ell_{q}^{n} \leqslant\|\mathcal{E}\|_{\ell_{p^{*}}^{n}} \ell_{p}^{n} .
$$

Consequently, since $\beta(p, q)=\beta(p, p)$ when $q \geqslant p$, the case $q \geqslant p$ follows from the case $q=p$.
So, suppose from now that $q \leqslant 2 \leqslant p$. If we denote

$$
r \stackrel{\text { def }}{=} \frac{q(p-2)}{p-q},
$$

with the convention $r=\infty$ if $q=p$, then $r \geqslant 1$ and

$$
\begin{equation*}
\frac{1}{q}=\frac{1-\theta}{r}+\frac{\theta}{2}, \quad \text { where } \quad \theta \stackrel{\operatorname{def}}{=} \frac{2}{p} \in[0,1] . \tag{462}
\end{equation*}
$$

Hence, by the Riesz-Thorin interpolation theorem [Rie27, Tho48] we have

$$
\|\mathcal{E}\|_{\ell_{p^{*}}^{n} \rightarrow \ell_{q}^{n}} \leqslant\|\mathcal{E}\|_{\ell_{1}^{n} \rightarrow \ell_{r}^{n}}^{1-\theta}\|\mathcal{E}\|_{\ell_{2}^{n} \rightarrow \ell_{2}^{n}}^{\theta}=\left(\max _{i \in\{1, \ldots, n\}}\left\|\mathcal{E} e_{i}\right\|_{\ell_{r}^{n}}\right)^{1-\theta}\|\mathcal{E}\|_{\ell_{2}^{n} \rightarrow \ell_{2}^{n}}^{\theta}=n^{\frac{1-\theta}{r}}\|\mathcal{E}\|_{\ell_{2}^{n} \rightarrow \ell_{2}^{n}}^{\theta} .
$$

By taking expectations of this inequality, we get that

$$
\begin{equation*}
\mathbb{E}\left[\|\mathcal{E}\|_{\ell_{p^{*}}^{n} \rightarrow \ell_{q}^{n}}\right] \leqslant n^{\frac{1-\theta}{r}} \mathbb{E}\left[\|\mathcal{E}\|_{\ell_{2}^{n} \rightarrow \ell_{2}^{n}}^{\theta}\right] \leqslant n^{\frac{1-\theta}{r}}\left(\mathbb{E}\left[\|\mathcal{E}\|_{\ell_{2}^{n} \rightarrow \ell_{2}^{n}}\right]\right)^{\theta} \lesssim n^{\frac{1-\theta}{r}+\frac{\theta}{2}}=n^{\frac{1}{q}}=n^{\beta(p, q)}, \tag{463}
\end{equation*}
$$

where the second step of (463) uses Jensen's inequality, the third step of (463) uses the classical fact that the expectation of the operator norm from $\ell_{2}^{n}$ to $\ell_{2}^{n}$ of an $n \times n$ matrix whose entries are i.i.d. symmetric Bernoulli random variables is $O(\sqrt{n})$ (this follows from [461), though it is older; see e.g. [BGN75]), the penultimate step of (463) uses (462), and the last step of (463) uses the definition of $\beta(p, q)$ in (457) while recalling that we are now treating the case $p \geqslant 2$ and $q \leqslant p$.

A substitution of Lemma 172 into the proof of Lemma 3.2 in [Sch82] yields the following asymptotic evaluations of the $n^{2}$-roots of volumes of the unit balls of injective and projective tensor products; the statement of [Sch82, Lemma 3.2] is identical, except that the constant factors depend on $p, q$, but that is due only to the dependence of the constants on $p, q$ in Lemma 2.3 in [Sch82], which Lemma 172 removes.

$$
\begin{equation*}
\operatorname{vol}_{n^{2}}\left(B_{\ell_{p}^{n} \check{\otimes} \ell_{q}^{n}}\right)^{\frac{1}{n^{2}}}=n^{-\beta(p, q)} \quad \text { and } \quad \operatorname{vol}_{n^{2}}\left(B_{\ell_{p}^{n} \stackrel{\otimes}{\otimes} \ell_{q}^{n}}\right)^{\frac{1}{n^{2}}}=n^{\beta\left(p^{*}, q^{*}\right)-2} \tag{464}
\end{equation*}
$$

Since $\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}$ belongs to the class of spaces in Remark 39 its Löwner ellipsoid is the minimal multiple of the standard Euclidean ball $B_{S_{2}^{n}}$ that superscribes the unit ball of $\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}$, namely

$$
\mathcal{L}_{\ell_{p}^{n} \hat{\otimes} \hat{\ell}_{q}^{n}}=R(n, p, q) B_{S_{2}^{n}}^{n},
$$

where, since $B_{\ell_{p}^{n}} \hat{\otimes} \ell_{q}^{n}$ is the convex hull of $B_{\ell_{p}^{n}} \otimes B_{\ell_{q}^{n}}$,

$$
\begin{equation*}
R(n, p, q)=\max _{\substack{x \in B_{n}^{n} \\ y \in B_{\ell_{q}^{n}}}}\|x \otimes y\|_{S_{2}^{n}}=\left(\max _{x \in B_{\ell_{p}^{n}}}\|x\|_{\ell_{2}^{n}}\right)\left(\max _{y \in B_{\ell_{q}^{n}}}\|y\|_{\ell_{2}^{n}}\right)=n^{\max \left\{\frac{1}{2}-\frac{1}{p}, 0\right\}+\max \left\{\frac{1}{2}-\frac{1}{q}, 0\right\}} . \tag{465}
\end{equation*}
$$

By combining (464) and (465) we get that

$$
\begin{align*}
& \operatorname{vr}\left(\ell_{p^{*}}^{n} \stackrel{\otimes}{\infty} \ell_{q^{*}}^{n}\right) \stackrel{[72]}{=} \operatorname{evr}\left(\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}\right)=R(n, p, q)\left(\frac{\operatorname{vol}_{n^{2}}\left(B_{S_{2}^{n}}\right)}{\operatorname{vol}_{n^{2}}\left(B_{\left.\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}\right)}\right.}\right)^{\frac{1}{n^{2}}} \tag{466}
\end{align*}
$$

A substitution of (466) into Theorem3 gives

$$
\operatorname{SEP}\left(\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}\right) \gtrsim \begin{cases}n^{\frac{3}{2}} & \text { if } \quad \max \{p, q\} \geqslant 2  \tag{467}\\ n^{\left.\left.1+\frac{1}{\max } \right\rvert\, p, q\right]} & \text { if }\end{cases}
$$

Furthermore, if Conjecture 11 holds for $\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}$, then (467) is sharp, namely (15) holds. Also, by Theorem 79 the left hand side of (467) is bounded from above by $O(\log n)$ times the right hand side of (467), thus implying the fifth bullet point of Corollary 4 .

Remark 173. The above results imply clustering statements (and impossibility thereof) for norms that have significance to algorithms and complexity theory. For example, the cut norm [FK99] on $\mathrm{M}_{n}(\mathbb{R})$ is $O(1)$-equivalent AN06] to the operator norm from $\ell_{\infty}^{n}$ to $\ell_{1}^{n}$. So, by (13) the separation modulus of the cut norm on $\mathrm{M}_{n}(\mathbb{R})$ is predicted to be bounded above and below by universal constant multiples of $n^{3 / 2}$, and by Theorem 79 we know that it is at least a universal constant multiple of $n^{3 / 2}$ and at most a universal constant multiple of $n^{3 / 2} \log n$. As another notable example, we proved that $\operatorname{SEP}\left(\ell_{\infty}^{n} \hat{\otimes} \ell_{\infty}^{n}\right) \gtrsim n^{3 / 2}$.

Moreover, if Conjecture 11 holds for $\ell_{\infty}^{n} \hat{\otimes} \ell_{\infty}^{n}$, then $\operatorname{SEP}\left(\ell_{\infty}^{n} \hat{\otimes} \ell_{\infty}^{n}\right)=n^{3 / 2}$ and by Theorem 79 we have $\operatorname{SEP}\left(\ell_{\infty}^{n} \hat{\otimes} \ell_{\infty}^{n}\right) \lesssim n^{3 / 2} \log n$. Grothendieck's inequality Gro53] implies that

$$
\begin{equation*}
\forall A \in \mathrm{M}_{n}(\mathbb{R}), \quad\|A\|_{\ell_{\infty}^{n} \hat{\otimes} \hat{\theta}_{\infty}^{n}}=\gamma_{2}^{1 \rightarrow \infty}(A), \tag{468}
\end{equation*}
$$

where $\gamma_{2}^{1 \rightarrow \infty}(A)$ is the factorization-through- $\ell_{2}$ norm (see [Pis86]) of $A$ as an operator from $\ell_{1}^{n}$ to $\ell_{\infty}^{n}$, i.e.,

$$
\gamma_{2}^{1 \rightarrow \infty}(A) \stackrel{\text { def }}{=} \min _{\substack{X, Y \in M_{n}(\mathbb{R}) \\ A=X Y}}\|X\|_{\ell_{2}^{n} \rightarrow \ell_{\infty}^{n}}\|Y\|_{\ell_{1}^{n} \rightarrow \ell_{2}^{n}} \min _{\substack{X, Y \in \mathrm{M}_{n}(\mathbb{R}) \\ A=X Y}} \max _{i \in\{1, \ldots, n\}}\left\|\operatorname{row}_{i}(X)\right\|_{\ell_{2}^{n}}\left\|\operatorname{column}_{j}(Y)\right\|_{\ell_{2}^{n}} .
$$

 $j$ 'th column of $M$, respectively. See [LMSS07] for the justification of (468), as well as the importance of the factorization norm $\gamma_{2}^{1 \rightarrow \infty}$ to complexity theory (see [MNT20, BLN21] for further algorithmic significance of factorization norms). Thanks to the above discussion, we know that

$$
n^{\frac{3}{2}} \lesssim \operatorname{SEP}\left(\mathrm{M}_{n}(\mathbb{R}), \gamma_{2}^{1 \rightarrow \infty}\right) \lesssim n^{\frac{3}{2}} \log n,
$$

and that $\operatorname{SEP}\left(\mathrm{M}_{n}(\mathbb{R}), \gamma_{2}^{1 \rightarrow \infty}\right)=n^{3 / 2}$ assuming Conjecture 11. To check that this does not follow from the previously known bounds (2), we need to know the asymptotic growth rate of the Banach-Mazur distance between $\ell_{\infty}^{n} \hat{\otimes} \ell_{\infty}^{n}$ and each of the spaces $\ell_{1}^{n^{2}}, \ell_{2}^{n^{2}}$. However, these Banach-Mazur distances do not appear in the literature. In response to our inquiry, Carsten Schütt answered this question, by showing that

$$
\begin{equation*}
d_{\mathrm{BM}}\left(\ell_{2}^{n^{2}}, \ell_{\infty}^{n} \hat{\otimes} \ell_{\infty}^{n}\right)=d_{\mathrm{BM}}\left(\ell_{1}^{n^{2}}, \ell_{\infty}^{n} \hat{\otimes} \ell_{\infty}^{n}\right)=n . \tag{469}
\end{equation*}
$$

More generally, Schütt succeeded to evaluate the asymptotic growth rate of the Banach-Mazur distance between $\ell_{p}^{n} \hat{\otimes} \ell_{q}^{n}$ and $\ell_{p}^{n} \dot{\otimes} \ell_{q}^{n}$ to each of $\ell_{1}^{n^{2}}, \ell_{2}^{n^{2}}$ for every $p, q \in[1, \infty]$ (this is a substantial matter that Schütt communicated to us privately and he will publish it elsewhere). Due to 469), an application of (2) only gives the bounds $n \lesssim \operatorname{SEP}\left(\ell_{\infty}^{n} \hat{\otimes} \ell_{\infty}^{n}\right) \lesssim n^{2}$, which hold for every $n^{2}$-dimensional normed space. More generally, Schütt's result shows that (13) and (15) do not follow from (2).

The volume computations of this section are only an indication of the available information. The literature contains many more volume estimates that could be substituted into Theorem 3 and Conjecture 6 to yield new results (and conjectures) on separation moduli of various spaces; examples of further pertinent results appear in [Sch82, Bal91, GJ97, GJN97, GJ99, GP99, GPS ${ }^{+} 17$, DP09, DV20, KP20, KPT20].

## 7. LOGARITHMIC WEAK ISOMORPHIC ISOPERIMETRY IN MINIMUM DUAL MEAN WIDTH POSITION

In this section we will prove the results that we stated in Section 1.6.3. We first claim that for every integer $n \geqslant 2$ and every $r>0$ we have

$$
\begin{equation*}
\operatorname{iq}\left(B_{\ell_{\infty}^{n}} \cap\left(r B_{\ell_{2}^{n}}\right)\right)=\operatorname{iq}\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right) \gtrsim\left(\min \{\sqrt{n}, r\}\left(1-\frac{1}{\max \left\{1, r^{2}\right\}}\right)^{\frac{n-1}{2}}+1\right) \sqrt{n} \tag{470}
\end{equation*}
$$

Observe that (470) implies (86). Furthermore, (470) implies the direction $\gtrsim$ in (87) because

$$
\begin{aligned}
\min _{r>0} \frac{\operatorname{iq}\left(B_{\ell_{\infty}^{n}} \cap\left(r B_{\ell_{2}^{n}}\right)\right)}{\sqrt{n}}\left(\frac{\operatorname{vol}_{n}\left(B_{\ell_{\infty}^{n}}\right)}{\operatorname{vol}_{n}\left(B_{\ell_{\infty}^{n}} \cap\left(r B_{\ell_{2}^{n}}\right)\right)}\right)^{\frac{1}{n}} & \geqslant \min _{r>0} \frac{\operatorname{iq}\left(B_{\ell_{\infty}^{n}} \cap\left(r B_{\ell_{2}^{n}}\right)\right)}{\sqrt{n}}\left(\frac{2^{n}}{\operatorname{vol}_{n}\left(r B_{\ell_{2}^{n}}\right)}\right)^{\frac{1}{n}} \\
& \gtrsim \min _{r>0}\left(\min \left\{\frac{\sqrt{n}}{r}, 1\right\}\left(1-\frac{1}{\max \left\{1, r^{2}\right\}}\right)^{\frac{n-1}{2}}+\frac{1}{r}\right) \sqrt{n}=\sqrt{\log n},
\end{aligned}
$$

where the penultimate step uses (470) and the final step is elementary calculus. Since the $K$-convexity constant of $\ell_{\infty}^{n}$ satisfies $K\left(\ell_{\infty}^{n}\right)=\sqrt{\log n}$ (see (Pis89, Chapter 2]), the matching upper bound in (87) will follow after we will prove (below) Proposition60. This will also show that Proposition60is sharp, though it would be worthwhile to find out if it is sharp even for some normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ for which $K(\mathbf{X})=\log n$; such a space exists by a remarkable (randomized) construction of Bourgain [Bou84].

To prove (470), note first that if $0<r \leqslant 1$, then $r B_{\ell_{2}^{n}} \subseteq[-1,1]^{n}$ and therefore

$$
\begin{equation*}
\forall 0<r \leqslant 1, \quad \operatorname{iq}\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right)=\operatorname{iq}\left(r B_{\ell_{2}^{n}}\right)=\sqrt{n} \tag{471}
\end{equation*}
$$

Similarly, note that if $r \geqslant \sqrt{n}$, then $r B_{\ell_{2}^{n}} \subseteq[-1,1]^{n}$ and therefore

$$
\begin{equation*}
\forall r \geqslant \sqrt{n}, \quad \operatorname{iq}\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right)=\operatorname{iq}\left([-1,1]^{n}\right)=n . \tag{472}
\end{equation*}
$$

Both (471) and (472) coincide with (470) in the respective ranges. The less trivial range of 470) is when $1<r<\sqrt{n}$, in which case the boundary of $[-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)$ contains the disjoint union of the intersection of $r B_{\ell_{2}^{n}}$ with the $2 n$ faces of $[-1,1]^{n}$, each of which is isometric to the following set.

$$
[-1,1]^{n-1} \cap\left(\sqrt{r^{2}-1} B_{\ell_{2}^{n-1}}\right)
$$

Together with the straightforward inclusion

$$
[-1,1]^{n-1} \cap\left(\sqrt{r^{2}-1} B_{\ell_{2}^{n-1}}\right) \supseteq \sqrt{1-\frac{1}{r^{2}}}\left([-1,1]^{n-1} \cap\left(r B_{\ell_{2}^{n-1}}\right)\right)
$$

the above observation implies that if $1<r<\sqrt{n}$, then

$$
\begin{align*}
\operatorname{vol}_{n-1}\left(\partial\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right)\right) & \geqslant 2 n\left(1-\frac{1}{r^{2}}\right)^{\frac{n-1}{2}} \operatorname{vol}_{n-1}\left([-1,1]^{n-1} \cap\left(r B_{\ell_{2}^{n-1}}\right)\right) \\
& =n\left(1-\frac{1}{r^{2}}\right)^{\frac{n-1}{2}} \operatorname{vol}_{n}\left(\left([-1,1]^{n-1} \cap\left(r B_{\ell_{2}^{n-1}}\right)\right) \times[-1,1]\right)  \tag{473}\\
& \geqslant n\left(1-\frac{1}{r^{2}}\right)^{\frac{n-1}{2}} \operatorname{vol}_{n}\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right),
\end{align*}
$$

where the final step (473) is a consequence of the straightforward inclusion

$$
\left([-1,1]^{n-1} \cap\left(r B_{\ell_{2}^{n-1}}\right)\right) \times[-1,1] \supseteq[-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)
$$

By combining (473) with the definition (11) of the isoperimetric quotient, we see that

$$
\begin{equation*}
\operatorname{iq}\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right) \geqslant \frac{n\left(1-\frac{1}{r^{2}}\right)^{\frac{n-1}{2}} \operatorname{vol}_{n}\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right)}{\operatorname{vol}_{n}\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right)^{\frac{n-1}{n}}}=n\left(1-\frac{1}{r^{2}}\right)^{\frac{n-1}{2}} \operatorname{vol}_{n}\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right)^{\frac{1}{n}} \tag{474}
\end{equation*}
$$

When $r \leqslant \sqrt{n}$ we have $[-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right) \supseteq[-r / \sqrt{n}, r / \sqrt{n}]^{n}$. In combination with 474), this implies that

$$
\forall 1<r<\sqrt{n}, \quad \operatorname{iq}\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right) \geqslant 2 r \sqrt{n}\left(1-\frac{1}{r^{2}}\right)^{\frac{n-1}{2}}
$$

As also $\operatorname{iq}\left([-1,1]^{n} \cap\left(r B_{\ell_{2}^{n}}\right)\right) \gtrsim \sqrt{n}$ by the isoperimetric theorem (12), this completes the proof of 470).
Passing to the proof of Proposition60, observe first that for every $r>0$ we have

$$
\begin{equation*}
\frac{\operatorname{vol}_{n}\left(B_{\mathbf{X}} \cap\left(r B_{\ell_{2}^{n}}\right)\right)}{\operatorname{vol}_{n}\left(r B_{\ell_{2}^{n}}\right.}=\frac{\operatorname{vol}_{n}\left(\left\{x \in r B_{\ell_{2}^{n}}:\|x\|_{\mathbf{X}} \leqslant 1\right\}\right)}{\operatorname{vol}_{n}\left(r B_{\ell_{2}^{n}}\right)} \geqslant 1-f_{r B_{\ell_{2}^{n}}}\|x\|_{\mathbf{X}} \mathrm{d} x=1-\frac{n r}{n+1} M(\mathbf{X}) \tag{475}
\end{equation*}
$$

where the penultimate step in (475) is Markov's inequality and the final step in (475) is integration in polar coordinates using the following standard notation for the mean of the norm on the Euclidean sphere:

$$
\begin{equation*}
M(\mathbf{X}) \stackrel{\text { def }}{=} f_{S^{n-1}}\|z\|_{\mathbf{X}} \mathrm{d} z \tag{476}
\end{equation*}
$$

We will also use the common notation $M^{*}(\mathbf{X}) \stackrel{\text { def }}{=} M\left(\mathbf{X}^{*}\right)$. By setting $r=1 /(2 M(\mathbf{X}))$ in (475) we get that

$$
\begin{equation*}
\operatorname{vol}_{n}\left(B_{\mathbf{X}} \cap\left(\frac{1}{2 M(\mathbf{X})} B_{\ell_{2}^{n}}\right)\right)^{\frac{1}{n}} \geqslant\left(\frac{1}{2} \operatorname{vol}_{n}\left(\frac{1}{2 M(\mathbf{X})} B_{\ell_{2}^{n}}\right)\right)^{\frac{1}{n}}=\frac{1}{M(\mathbf{X}) \sqrt{n}} \tag{477}
\end{equation*}
$$

This simple consideration gives the following general elementary lemma.
Lemma 174. Let $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$ be a normed space. For $r=1 /(2 M(\mathbf{X}))$ and $L=B_{\mathbf{X}} \cap\left(r B_{\ell_{2}^{n}}\right)$ we have

$$
\begin{equation*}
\operatorname{vol}_{n}(L)^{\frac{1}{n}} \gtrsim \frac{1}{M(\mathbf{X}) \sqrt{n}} \quad \text { and } \quad \frac{\operatorname{MaxProj}(L)}{\operatorname{vol}_{n}(L)^{\frac{n-1}{n}}} \lesssim 1 \tag{478}
\end{equation*}
$$

Proof. The first inequality in (478) follows from (477). For the second inequality in (478), observe that $\operatorname{Proj}_{z^{\perp}}(L) \subseteq \operatorname{Proj}_{z^{\perp}}\left(r B_{\ell_{2}^{n}}\right)$ for every $z \in S^{n-1}$, since $L \subseteq r B_{\ell_{2}^{n}}$. Consequently,

$$
\begin{equation*}
\operatorname{MaxProj}(L) \leqslant \operatorname{MaxProj}\left(r B_{\ell_{2}^{n}}\right)=r^{n-1} \operatorname{vol}_{n-1}\left(B_{\ell_{2}^{n-1}}\right) \asymp \operatorname{vol}_{n}\left(r B_{\ell_{2}^{n}}\right)^{\frac{n-1}{n}} \lesssim \operatorname{vol}_{n}(L)^{\frac{n-1}{n}} \tag{479}
\end{equation*}
$$

where the penultimate step of (479) is a standard computation using Stirling's formula and the final step of (479) uses the first inequality in (477).

By (54), the second inequality in (478) implies that $\mathrm{iq}(L) \lesssim \sqrt{n}$. Hence, in order to use Lemma 174 in the context of Conjecture 10 it would be beneficial to choose $S \in S L_{n}(\mathbb{R})$ for which $M(S X)$ is small. So, fix $\delta>0$ and suppose that $\delta M(S \mathbf{X}) \leqslant \min _{T \in \mathrm{SL}_{n}(\mathbb{R})} M(T \mathbf{X})$. By compactness, this holds for some $S \in \mathrm{SL}_{n}(\mathbb{R})$ with $\delta=1$, in which case the polar of $S B_{\mathbf{X}}$ is in minimum mean width position and we will say that $S \mathbf{X}$ is in minimum dual mean width position (the terminology that is used in [GMR00] is that $S B_{\mathbf{X}}$ has minimal $M)$. By [GM00], the matrix in $\mathrm{SL}_{n}(\mathbb{R})$ at which $\min _{T \in S L_{n}(\mathbb{R})} M(T \mathbf{X})$ is attained is unique up to orthogonal transformations. We allow the flexibility of working with some universal constant $0<\delta<1$ rather than considering only the minimum dual mean width position since this will encompass other commonly used positions, such as the $\ell$-position (see BGVV14, Section 1.11]). By [GM00], $\mathbf{X}$ is in minimum dual mean width position if and only if the measure $\mathrm{d} v_{\mathbf{X}}(z)=\|z\|_{\mathbf{X}} \mathrm{d} z$ on $S^{n-1}$ is isotropic. Since $v_{\mathbf{X}}$ is evidently Isom ( $\mathbf{X}$ )-invariant, by (70) if $\mathbf{X}$ is canonically positioned, then it is in minimum dual mean width position.

Let $\gamma$ be the standard Gaussian measure on $\mathbb{R}$, i.e., its density is $u \mapsto \exp \left(-u^{2} / 2\right) / \sqrt{2 \pi}$. The (Gaussian) $K$-convexity constant $K(\mathbf{X})$ of $\mathbf{X}$ is defined MP76] to be the infimum over those $K>0$ that satisfy

$$
\left(\int_{\mathbb{R}^{\aleph_{0}}}\left\|\sum_{i=1}^{\infty} \mathrm{g}_{i}^{\prime} \int_{\mathbb{R}^{\aleph_{0}}} \mathrm{~g}_{i} f(\mathrm{~g}) \mathrm{d} \gamma^{\otimes \aleph_{0}}(\mathrm{~g})\right\|_{\mathbf{X}}^{2} \mathrm{~d} \gamma^{\otimes \aleph_{0}}\left(\mathrm{~g}^{\prime}\right)\right)^{\frac{1}{2}} \leqslant K\left(\int_{\mathbb{R}^{\aleph_{0}}}\|f(\mathrm{~g})\|_{\mathbf{X}}^{2} \mathrm{~d} \gamma^{\otimes \aleph_{0}}(\mathrm{~g})\right)^{\frac{1}{2}}
$$

for every measurable $f: \mathbb{R}^{\aleph_{0}} \rightarrow \mathbf{X}$ with $\int_{\mathbb{R}^{\aleph_{0}}}\|f(\mathrm{~g})\|_{\mathrm{X}}^{2} \mathrm{~d} \gamma^{\otimes \aleph_{0}}(\mathrm{~g})<\infty$. By FTJ79 there is $T \in \mathrm{SL}_{n}(\mathbb{R})$ such that $M(T \mathbf{X}) M^{*}(T \mathbf{X}) \leqslant K(\mathbf{X})$. By the above assumption $\delta M(S \mathbf{X}) \leqslant M(T \mathbf{X})$, so $\delta M(S \mathbf{X}) \leqslant K(\mathbf{X}) / M^{*}(T \mathbf{X})$. Next, $M(\mathbf{X}) \geqslant\left(\operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right) / \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)\right)^{1 / n}$; see e.g. MP89, Section 2] and HN19, Lemma 30] for two derivations of this well-known volumetric lower bound on $M(\mathbf{X})$. Applying this lower bound to the dual of $T \mathbf{X}$, we get $M^{*}(T \mathbf{X}) \geqslant\left(\operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right) / \operatorname{vol}_{n}\left(B_{\mathbf{X}^{*}}\right)\right)^{1 / n}$. The Blaschke-Santaló inequality Bla17, San49] states that

$$
\frac{\operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}^{*}}\right)} \geqslant \frac{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\ell_{2}^{n}}\right)}
$$

so we conclude that $\delta M(S \mathbf{X}) \sqrt{n} \lesssim K(\mathbf{X}) / \sqrt[n]{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)}$. A substitution of this into Lemma 174 gives:
Proposition 175. Fix $0<\delta \leqslant 1$ and a normed space $\mathbf{X}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{X}}\right)$. Suppose that $S \in \mathrm{SL}_{n}(\mathbb{R})$ satisfies $\delta M(S \mathbf{X}) \leqslant \min _{T \in \mathrm{SL}_{n}(\mathbb{R})} M(T \mathbf{X})$. Then, denoting $r=1 /(2 M(S \mathbf{X}))$ we have

$$
\operatorname{vol}_{n}\left(\left(S B_{\mathbf{X}}\right) \cap\left(r B_{\ell_{2}^{n}}\right)\right)^{\frac{1}{n}} \gtrsim \frac{\delta}{K(\mathbf{X})} \operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{\frac{1}{n}} \quad \text { and } \quad \operatorname{MaxProj}\left(\left(S B_{\mathbf{X}}\right) \cap\left(r B_{\ell_{2}^{n}}\right)\right) \lesssim \operatorname{vol}_{n}\left(\left(S B_{\mathbf{X}}\right) \cap\left(r B_{\ell_{2}^{n}}\right)\right)^{\frac{n-1}{n}}
$$

Furthermore, if $\mathbf{X}$ is canonically positioned, then this holds when $S$ is the identity matrix and $\delta=1$.
By (54), Proposition 175 implies Proposition 60 , with the additional information that the conclusion of Proposition 60 holds with $S$ the identity matrix if $\mathbf{X}$ is in minimum dual mean width position, in which case we obtain an upper bound on $\operatorname{MaxProj}(L)$. Hence, by the reasoning in Section 1.6, if $\mathbf{X}$ is in minimum dual mean width position, then

$$
\operatorname{SEP}(\mathbf{X}) \lesssim K(\mathbf{X}) \frac{\operatorname{diam}_{\ell_{2}^{n}}\left(B_{\mathbf{X}}\right)}{\operatorname{vol}_{n}\left(B_{\mathbf{X}}\right)^{\frac{1}{n}}}
$$

## References

[AAR99] G. E. Andrews, R. Askey, and R. Roy. Special functions, volume 71 of Encyclopedia of Mathematics and its Applications. Cambridge University Press, Cambridge, 1999. ISBN 0-521-62321-9; 0-521-78988-5. doi:10.1017/ CBO9781107325937.
[ABN11] I. Abraham, Y. Bartal, and O. Neiman. Advances in metric embedding theory. Adv. Math., 228(6):3026-3126, 2011.
[AC09] F. Alter and V. Caselles. Uniqueness of the Cheeger set of a convex body. Nonlinear Anal., 70(1):32-44, 2009.
[ACC05] F. Alter, V. Caselles, and A. Chambolle. A characterization of convex calibrable sets in $\mathbb{R}^{N}$. Math. Ann., 332(2):329-366, 2005.
[Ach03] D. Achlioptas. Database-friendly random projections: Johnson-Lindenstrauss with binary coins. J. Comput. System Sci., 66(4):671-687, 2003. Special issue on PODS 2001 (Santa Barbara, CA).
[AE56] R. F. Arens and J. Eells, Jr. On embedding uniform and topological spaces. Pacific J. Math., 6:397-403, 1956.
[AFP00] L. Ambrosio, N. Fusco, and D. Pallara. Functions of bounded variation and free discontinuity problems. Oxford Mathematical Monographs. The Clarendon Press, Oxford University Press, New York, 2000. ISBN 0-19-850245-1.
[AGS08] L. Ambrosio, N. Gigli, and G. Savaré. Gradient flows in metric spaces and in the space of probability measures. Lectures in Mathematics ETH Zürich. Birkhäuser Verlag, Basel, second edition, 2008. ISBN 978-3-7643-8721-1.
[AGZ10] G. W. Anderson, A. Guionnet, and O. Zeitouni. An introduction to random matrices, volume 118 of Cambridge Studies in Advanced Mathematics. Cambridge University Press, Cambridge, 2010. ISBN 978-0-521-19452-5.
[AI06] A. Andoni and P. Indyk. Near-optimal hashing algorithms for approximate nearest neighbor in high dimensions. In 47th Annual IEEE Symposium on Foundations of Computer Science (FOCS 2006), 21-24 October 2006, Berkeley, California, USA, Proceedings, pages 459-468. IEEE Computer Society, 2006.
[AK06] F. Albiac and N. J. Kalton. Topics in Banach space theory, volume 233 of Graduate Texts in Mathematics. Springer, New York, 2006. ISBN 978-0387-28141-4; 0-387-28141-X.
[AKPW91] N. Alon, R. M. Karp, D. Peleg, and D. B. West. A graph-theoretic game and its application to the k-server problem (extended abstract). In L. A. McGeoch and D. D. Sleator, editors, On-Line Algorithms, Proceedings of a DIMACS Workshop, New Brunswick, New Jersey, USA, February 11-13, 1991, volume 7 of DIMACS Series in Discrete Mathematics and Theoretical Computer Science, pages 1-10. DIMACS/AMS, 1991.
[AM83] N. Alon and V. D. Milman. Embedding of $l_{\infty}^{k}$ in finite-dimensional Banach spaces. Israel J. Math., 45(4):265-280, 1983.
[AN06] N. Alon and A. Naor. Approximating the cut-norm via Grothendieck's inequality. SIAM J. Comput., 35(4):787-803, 2006.
[AP90] B. Awerbuch and D. Peleg. Sparse partitions (extended abstract). In 31st Annual Symposium on Foundations of Computer Science, St. Louis, Missouri, USA, October 22-24, 1990, Volume II, pages 503-513. IEEE Computer Society, 1990. doi:10.1109/FSCS.1990.89571.
[AP20] L. Ambrosio and D. Puglisi. Linear extension operators between spaces of Lipschitz maps and optimal transport. J. Reine Angew. Math., 764:1-21, 2020.
[ATTJ05] R. Anisca, A. Tcaciuc, and N. Tomczak-Jaegermann. Structure of normed spaces with extremal distance to the Euclidean space. Houston J. Math., 31(1):267-283, 2005.
[Aus11] T. Austin. A CAT(0)-valued pointwise ergodic theorem. J. Topol. Anal., 3(2):145-152, 2011.
[Bal89] K. Ball. Volumes of sections of cubes and related problems. In Geometric aspects of functional analysis (1987-88), volume 1376 of Lecture Notes in Math., pages 251-260. Springer, Berlin, 1989. doi:10.1007/BFb0090058.
[Bal91a] K. Ball. Normed spaces with a weak-Gordon-Lewis property. In Functional analysis (Austin, TX, 1987/1989), volume 1470 of Lecture Notes in Math., pages 36-47. Springer, Berlin, 1991. doi:10.1007/BFb0090210.
[Bal91b] K. Ball. Shadows of convex bodies. Trans. Amer. Math. Soc., 327(2):891-901, 1991.
[Bal91c] K. Ball. Volume ratios and a reverse isoperimetric inequality. J. London Math. Soc. (2), 44(2):351-359, 1991.
[Bal92] K. Ball. Markov chains, Riesz transforms and Lipschitz maps. Geom. Funct. Anal., 2(2):137-172, 1992.
[Bal01] K. Ball. Convex geometry and functional analysis. In Handbook of the geometry of Banach spaces, Vol. I, pages 161194. North-Holland, Amsterdam, 2001. doi:10.1016/S1874-5849(01)80006-1.
[Ban93] S. Banach. Théorie des opérations linéaires. Éditions Jacques Gabay, Sceaux, 1993. ISBN 2-87647-148-5. Reprint of the 1932 original.
[Bar96] Y. Bartal. Probabilistic approximation of metric spaces and its algorithmic applications. In 37th Annual Symposium on Foundations of Computer Science (Burlington, VT, 1996), pages 184-193. IEEE Comput. Soc. Press, Los Alamitos, CA, 1996.
[Bar99] Y. Bartal. On approximating arbitrary metrices by tree metrics. In STOC '98 (Dallas, TX), pages 161-168. ACM, New York, 1999.
[Bas18] G. Basso. Fixed point theorems for metric spaces with a conical geodesic bicombing. Ergodic Theory Dynam. Systems, 38(5):1642-1657, 2018.
[BB05] A. Brudnyĭ and Y. Brudny̌̌. Simultaneous extensions of Lipschitz functions. Uspekhi Mat. Nauk, 60(6(366)):53-72, 2005.
[BB06] A. Brudnyi and Y. Brudnyi. Extension of Lipschitz functions defined on metric subspaces of homogeneous type. Rev. Mat. Complut., 19(2):347-359, 2006.
[BB07a] A. Brudnyi and Y. Brudnyi. Linear and nonlinear extensions of Lipschitz functions from subsets of metric spaces. Algebra i Analiz, 19(3):106-118, 2007.
[BB07b] A. Brudnyi and Y. Brudnyi. Metric spaces with linear extensions preserving Lipschitz condition. Amer. J. Math., 129(1):217-314, 2007.
[BB12] A. Brudnyi and Y. Brudnyi. Methods of geometric analysis in extension and trace problems. Volume 2, volume 103 of Monographs in Mathematics. Birkhäuser/Springer Basel AG, Basel, 2012.
[BC05] B. Brinkman and M. Charikar. On the impossibility of dimension reduction in $l_{1}$. J. ACM, 52(5):766-788 (electronic), 2005.
[Ben77] G. Bennett. Schur multipliers. Duke Math. J., 44(3):603-639, 1977.
[BF87] I. Bárány and Z. Füredi. Computing the volume is difficult. Discrete Comput. Geom., 2(4):319-326, 1987.
[BGMN05] F. Barthe, O. Guédon, S. Mendelson, and A. Naor. A probabilistic approach to the geometry of the $l_{p}^{n}$-ball. Ann. Probab., 33(2):480-513, 2005.
[BGN75] G. Bennett, V. Goodman, and C. M. Newman. Norms of random matrices. Pacific J. Math., 59(2):359-365, 1975.
[BGVV14] S. Brazitikos, A. Giannopoulos, P. Valettas, and B.-H. Vritsiou. Geometry of isotropic convex bodies, volume 196 of Mathematical Surveys and Monographs. American Mathematical Society, Providence, RI, 2014. ISBN 978-1-4704-1456-6. doi:10.1090/surv/196.
[BH99] M. R. Bridson and A. Haefliger. Metric spaces of non-positive curvature, volume 319 of Grundlehren der Mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences]. Springer-Verlag, Berlin, 1999.
[Bha97] R. Bhatia. Matrix analysis, volume 169 of Graduate Texts in Mathematics. Springer-Verlag, New York, 1997. ISBN 0-387-94846-5. doi:10.1007/978-1-4612-0653-8.
[BK81] J. A. Brudnyĭ and N. J. Krugljak. Functors of real interpolation. Dokl. Akad. Nauk SSSR, 256(1):14-17, 1981.
[BL88] J. Bourgain and J. Lindenstrauss. Projection bodies. In Geometric aspects of functional analysis (1986/87), volume 1317 of Lecture Notes in Math., pages 250-270. Springer, Berlin, 1988. doi:10.1007/BFb0081746.
[BL00] Y. Benyamini and J. Lindenstrauss. Geometric nonlinear functional analysis. Vol. 1, volume 48 of American Mathematical Society Colloquium Publications. American Mathematical Society, Providence, RI, 2000.
[Bla17] W. Blaschke. Über affine Geometrie VII: Neue Extremeigenschaften von Ellipse und Ellipsoid. Leipz. Ber. 69, 306-318 (1917)., 1917.
[BLM89] J. Bourgain, J. Lindenstrauss, and V. Milman. Approximation of zonoids by zonotopes. Acta Math., 162(1-2):73-141, 1989.
[BLN21] V. Bhattiprolu, E. Lee, and A. Naor. A framework for quadratic form maximization over convex sets through nonconvex relaxations. In S. Khuller and V. V. Williams, editors, STOC '21: 53rd Annual ACM SIGACT Symposium on Theory of Computing, Virtual Event, Italy, June 21-25, 2021, pages 870-881. ACM, 2021. doi:10.1145/3406325.3451128.
[BM85] J. Bourgain and V. D. Milman. Sections euclidiennes et volume des corps symétriques convexes dans $\mathbf{R}^{n}$. C. R. Acad. Sci. Paris Sér. I Math., 300(13):435-438, 1985.
[BM87] J. Bourgain and V. D. Milman. New volume ratio properties for convex symmetric bodies in $\mathbf{R}^{n}$. Invent. Math., 88(2):319-340, 1987.
[BN02] F. Barthe and A. Naor. Hyperplane projections of the unit ball of $\ell_{p}^{n}$. Discrete Comput. Geom., 27(2):215-226, 2002.
[Boc33] S. Bochner. Integration von Funktionen, deren Werte die Elemente eines Vektorraumes sind. Fundam. Math., 20:262-276, 1933.
[Bol69] E. D. Bolker. A class of convex bodies. Trans. Amer. Math. Soc., 145:323-345, 1969.
[Bol08] F. Bolley. Separability and completeness for the Wasserstein distance. In Séminaire de probabilités XLI, volume 1934 of Lecture Notes in Math., pages 371-377. Springer, Berlin, 2008. doi:10.1007/978-3-540-77913-1_17.
[Bou82] J. Bourgain. A remark on finite-dimensional $P_{\lambda}$-spaces. Studia Math., 72(3):285-289, 1982.
[Bou84] J. Bourgain. On martingales transforms in finite-dimensional lattices with an appendix on the $K$-convexity constant. Math. Nachr., 119:41-53, 1984.
[Bou87] J. Bourgain. On dimension free maximal inequalities for convex symmetric bodies in $\mathbf{R}^{n}$. In Geometrical aspects of functional analysis (1985/86), volume 1267 of Lecture Notes in Math., pages 168-176. Springer, Berlin, 1987.
[BP90] K. Ball and A. Pajor. Convex bodies with few faces. Proc. Amer. Math. Soc., 110(1):225-231, 1990.
[Bra20] L. Brasco. On principal frequencies and isoperimetric ratios in convex sets. Ann. Fac. Sci. Toulouse Math. (6), 29(4):977-1005, 2020.
[Bro41] R. L. Brooks. On colouring the nodes of a network. Proc. Cambridge Philos. Soc., 37:194-197, 1941.
[BS88] J. Bourgain and S. J. Szarek. The Banach-Mazur distance to the cube and the Dvoretzky-Rogers factorization. Israel J. Math., 62(2):169-180, 1988.
[BS02] Y. Brudnyi and P. Shvartsman. Stability of the Lipschitz extension property under metric transforms. Geom. Funct. Anal., 12(1):73-79, 2002.
[BT87] J. Bourgain and L. Tzafriri. Invertibility of "large" submatrices with applications to the geometry of Banach spaces and harmonic analysis. Israel J. Math., 57(2):137-224, 1987.
[Bus82] P. Buser. A note on the isoperimetric constant. Ann. Sci. École Norm. Sup. (4), 15(2):213-230, 1982.
[But72] G. J. Butler. Simultaneous packing and covering in euclidean space. Proc. London Math. Soc. (3), 25:721-735, 1972.
[Car85] B. Carl. Inequalities of Bernstein-Jackson-type and the degree of compactness of operators in Banach spaces. Ann. Inst. Fourier (Grenoble), 35(3):79-118, 1985.
[CCG $\left.{ }^{+} 98\right]$ M. Charikar, C. Chekuri, A. Goel, S. Guha, and S. A. Plotkin. Approximating a finite metric by a small number of tree metrics. In 39th Annual Symposium on Foundations of Computer Science, FOCS '98, November 8-11, 1998, Palo Alto, California, USA, pages 379-388. IEEE Computer Society, 1998.
[CCN07] V. Caselles, A. Chambolle, and M. Novaga. Uniqueness of the Cheeger set of a convex body. Pacific J. Math., 232(1):7790, 2007.
[CF86] G. D. Chakerian and P. Filliman. The measures of the projections of a cube. Studia Sci. Math. Hungar., 21(1-2):103110, 1986.
[CGFS86] F. R. K. Chung, R. L. Graham, P. Frankl, and J. B. Shearer. Some intersection theorems for ordered sets and graphs. J. Combin. Theory Ser. A, 43(1):23-37, 1986.
[CH53] R. Courant and D. Hilbert. Methods of mathematical physics. Vol. I. Interscience Publishers, Inc., New York, N.Y., 1953.
[Cha84] I. Chavel. Eigenvalues in Riemannian geometry, volume 115 of Pure and Applied Mathematics. Academic Press, Inc., Orlando, FL, 1984. ISBN 0-12-170640-0. Including a chapter by Burton Randol, With an appendix by Jozef Dodziuk.
[Che70] J. Cheeger. A lower bound for the smallest eigenvalue of the Laplacian. In Problems in analysis (Papers dedicated to Salomon Bochner, 1969), pages 195-199. Princeton Univ. Press, Princeton, N. J., 1970.
[Che78] S. Chevet. Séries de variables aléatoires gaussiennes à valeurs dans $E \hat{\otimes}_{\mathcal{E}} F$. Application aux produits d'espaces de Wiener abstraits. In Séminaire sur la Géométrie des Espaces de Banach (1977-1978), pages Exp. No. 19, 15. École Polytech., Palaiseau, 1978.
[CKR05] G. Calinescu, H. Karloff, and Y. Rabani. Approximation algorithms for the 0 -extension problem. SIAM J. Comput., 34(2):358-372, 2004/05.
[CP88] B. Carl and A. Pajor. Gel'fand numbers of operators with values in a Hilbert space. Invent. Math., 94(3):479-504, 1988.
[Cwi84] M. Cwikel. $K$-divisibility of the $K$-functional and Calderón couples. Ark. Mat., 22(1):39-62, 1984.
[Day47] M. M. Day. Polygons circumscribed about closed convex curves. Trans. Amer. Math. Soc., 62:315-319, 1947.
[DFS08] J. Diestel, J. H. Fourie, and J. Swart. The metric theory of tensor products. American Mathematical Society, Providence, RI, 2008. ISBN 978-0-8218-4440-3. doi:10.1090/mbk/052. Grothendieck's résumé revisited.
[DL15] D. Descombes and U. Lang. Convex geodesic bicombings and hyperbolicity. Geom. Dedicata, 177:367-384, 2015.
[DM06] A. Defant and C. Michels. Norms of tensor product identities. Note Mat., 25(1):129-166, 2005/06.
[DP09] A. Defant and C. Prengel. Volume estimates in spaces of homogeneous polynomials. Math. Z., 261(4):909-932, 2009.
[DR50] A. Dvoretzky and C. A. Rogers. Absolute and unconditional convergence in normed linear spaces. Proc. Nat. Acad. Sci. U. S. A., 36:192-197, 1950.
[Dur19] R. Durrett. Probability-theory and examples, volume 49 of Cambridge Series in Statistical and Probabilistic Mathematics. Cambridge University Press, Cambridge, 2019. ISBN 978-1-108-47368-2. doi:10.1017/9781108591034. Fifth edition of [ MR1068527].
[DV20] A. Doležalová and J. Vybíral. On the volume of unit balls of finite-dimensional Lorentz spaces. J. Approx. Theory, 255:105407, 20, 2020.
[Dvo61] A. Dvoretzky. Some results on convex bodies and Banach spaces. In Proc. Internat. Sympos. Linear Spaces (Jerusalem, 1960), pages 123-160. Jerusalem Academic Press, Jerusalem; Pergamon, Oxford, 1961.
[Elt83] J. Elton. Sign-embeddings of $l_{1}^{n}$. Trans. Amer. Math. Soc., 279(1):113-124, 1983.
[Enf70] P. Enflo. Uniform structures and square roots in topological groups. Part II. Israel J. Math. 8 (1970), 230-252; ibid., 8:253-272, 1970.
[ESH99] A. Es-Sahib and H. Heinich. Barycentre canonique pour un espace métrique à courbure négative. In Séminaire de Probabilités, XXXIII, volume 1709 of Lecture Notes in Math., pages 355-370. Springer, Berlin, 1999.
[Fab23] G. Faber. Beweis, daß unter allen homogenen Membranen von gleicher Fläche und gleicher Spannung die kreisförmige den tiefsten Grundton gibt. Münch. Ber. 1923, 169-172 (1923)., 1923.
[Fig77] Figiel, T., 1977. Review of KM73], Math. Rev. 53, \# 3649.
[FJ80] T. Figiel and W. B. Johnson. Large subspaces of $l_{\infty}^{n}$ and estimates of the Gordon-Lewis constant. IsraelJ. Math., 37(1-2):92-112, 1980.
[FK99] A. Frieze and R. Kannan. Quick approximation to matrices and applications. Combinatorica, 19(2):175-220, 1999.
[FLM77] T. Figiel, J. Lindenstrauss, and V. D. Milman. The dimension of almost spherical sections of convex bodies. Acta Math., 139(1-2):53-94, 1977.
[FRT04] J. Fakcharoenphol, S. Rao, and K. Talwar. A tight bound on approximating arbitrary metrics by tree metrics. J. Comput. System Sci., 69(3):485-497, 2004.
[FTJ79] T. Figiel and N. Tomczak-Jaegermann. Projections onto Hilbertian subspaces of Banach spaces. Israel J. Math., 33(2):155-171, 1979.
[Gar06] R. J. Gardner. Geometric tomography, volume 58 of Encyclopedia of Mathematics and its Applications. Cambridge University Press, New York, second edition, 2006. ISBN 0-521; 0-521-68493-5. doi:10.1017/CBO9781107341029.
[GG71] D. J. H. Garling and Y. Gordon. Relations between some constants associated with finite dimensional Banach spaces. Israel J. Math., 9:346-361, 1971.
[Gia95] A. A. Giannopoulos. A note on the Banach-Mazur distance to the cube. In Geometric aspects of functional analysis (Israel, 1992-1994), volume 77 of Oper. Theory Adv. Appl., pages 67-73. Birkhäuser, Basel, 1995.
[Gia96] A. A. Giannopoulos. A proportional Dvoretzky-Rogers factorization result. Proc. Amer. Math. Soc., 124(1):233-241, 1996.
[GJ97] Y. Gordon and M. Junge. Volume formulas in $L_{p}$-spaces. Positivity, 1(1):7-43, 1997.
[GJ99] Y. Gordon and M. Junge. Volume ratios in $L_{p}$-spaces. Studia Math., 136(2):147-182, 1999.
[GJN97] Y. Gordon, M. Junge, and N. J. Nielsen. The relations between volume ratios and new concepts of GL constants. Positivity, 1(4):359-379, 1997.
[GK03] G. Godefroy and N. J. Kalton. Lipschitz-free Banach spaces. Studia Math., 159(1):121-141, 2003. Dedicated to Professor Aleksander Pełczyński on the occasion of his 70th birthday.
[GKL03] A. Gupta, R. Krauthgamer, and J. R. Lee. Bounded geometries, fractals, and low-distortion embeddings. In 44th Symposium on Foundations of Computer Science (FOCS 2003), 11-14 October 2003, Cambridge, MA, USA, Proceedings, pages 534-543. IEEE Computer Society, 2003.
[GKM66] V. I. Gurariĭ, M. Ĭ. Kadec', and V. I. Macaev. Distances between finite-dimensional analogs of the $L_{p}$-spaces. Mat. Sb. (N.S.), 70 (112):481-489, 1966.
[GL79] Y. Gordon and R. Loewy. Uniqueness of ( $\Delta$ ) bases and isometries of Banach spaces. Math. Ann., 241(2):159-180, 1979.
[Glu81] E. D. Gluskin. The diameter of the Minkowski compactum is roughly equal to n. Funktsional. Anal. i Prilozhen., 15(1):72-73, 1981.
[Glu88] E. D. Gluskin. Extremal properties of orthogonal parallelepipeds and their applications to the geometry of Banach spaces. Mat. Sb. (N.S.), 136(178)(1):85-96, 1988.
[GM87] M. Gromov and V. D. Milman. Generalization of the spherical isoperimetric inequality to uniformly convex Banach spaces. Compositio Math., 62(3):263-282, 1987.
[GM00] A. A. Giannopoulos and V. D. Milman. Extremal problems and isotropic positions of convex bodies. Israel J. Math., 117:29-60, 2000.
[GMR00] A. A. Giannopoulos, V. D. Milman, and M. Rudelson. Convex bodies with minimal mean width. In Geometric aspects of functional analysis, volume 1745 of Lecture Notes in Math., pages 81-93. Springer, Berlin, 2000. doi: 10.1007/BFb0107209.
[GNS12] O. Giladi, A. Naor, and G. Schechtman. Bourgain's discretization theorem. Ann. Fac. Sci. Toulouse Math. (6), 21(4):817-837, 2012. Available athttps://arxiv.org/abs/1110.5368
[God15] G. Godefroy. A survey on Lipschitz-free Banach spaces. Comment. Math., 55(2):89-118, 2015.
[GP99] A. Giannopoulos and M. Papadimitrakis. Isotropic surface area measures. Mathematika, 46(1):1-13, 1999.
[GPS ${ }^{+}$17] O. Giladi, J. Prochno, C. Schütt, N. Tomczak-Jaegermann, and E. Werner. On the geometry of projective tensor products. J. Funct. Anal., 273(2):471-495, 2017.
[Gro53] A. Grothendieck. Résumé de la théorie métrique des produits tensoriels topologiques. Bol. Soc. Mat. São Paulo, 8:179, 1953.
[Gro03] M. Gromov. Random walk in random groups. Geom. Funct. Anal., 13(1):73-146, 2003.
[Grü60] B. Grünbaum. Projection constants. Trans. Amer. Math. Soc., 95:451-465, 1960.
[Had57] H. Hadwiger. Vorlesungen über Inhalt, Oberfläche und Isoperimetrie. Springer-Verlag, Berlin-Göttingen-Heidelberg, 1957.
[Hen06] A. Henrot. Extremum problems for eigenvalues of elliptic operators. Frontiers in Mathematics. Birkhäuser Verlag, Basel, 2006. ISBN 978-3-7643-7705-2; 3-7643-7705-4.
[Hen12] M. Henk. Löwner-John ellipsoids. Doc. Math., (Extra vol.: Optimization stories):95-106, 2012.
[HL34] G. H. Hardy and J. E. Littlewood. Bilinear forms bounded in spaces (p, q). Quart. J. Math. Oxford, 5(1):241-254, 1934.
[HLN16] T. Hytönen, S. Li, and A. Naor. Quantitative affine approximation for UMD targets. Discrete Anal., pages Paper No. 6, 37, 2016.
[HN19] T. Hytönen and A. Naor. Heat flow and quantitative differentiation. J. Eur. Math. Soc. (JEMS), 21(11):3415-3466, 2019.
[HP18] A. Henrot and M. Pierre. Shape variation and optimization, volume 28 of EMS Tracts in Mathematics. European Mathematical Society (EMS), Zürich, 2018. ISBN 978-3-03719-178-1. doi:10.4171/178.
[Jac68] M. Q. Jacobs. Measurable multivalued mappings and Lusin's theorem. Trans. Amer. Math. Soc., 134:471-481, 1968.
[Jam78] R. C. James. Nonreflexive spaces of type 2. Israel J. Math., 30(1-2):1-13, 1978.
[Jec03] T. Jech. Set theory. Springer Monographs in Mathematics. Springer-Verlag, Berlin, 2003. The third millennium edition, revised and expanded.
[JL84] W. B. Johnson and J. Lindenstrauss. Extensions of Lipschitz mappings into a Hilbert space. In Conference in modern analysis and probability (New Haven, Conn., 1982), volume 26 of Contemp. Math., pages 189-206. Amer. Math. Soc., Providence, RI, 1984.
[JL01] W. B. Johnson and J. Lindenstrauss. Basic concepts in the geometry of Banach spaces. In Handbook of the geometry of Banach spaces, Vol. I, pages 1-84. North-Holland, Amsterdam, 2001. doi:10.1016/S1874-5849(01)80003-6.
[JLS86] W. B. Johnson, J. Lindenstrauss, and G. Schechtman. Extensions of Lipschitz maps into Banach spaces. Israel J. Math., 54(2):129-138, 1986.
[JN10] W. B. Johnson and A. Naor. The Johnson-Lindenstrauss lemma almost characterizes Hilbert space, but not quite. Discrete Comput. Geom., 43(3):542-553, 2010.
[Joh48] F. John. Extremum problems with inequalities as subsidiary conditions. In Studies and Essays Presented to R. Courant on his 60th Birthday, January 8, 1948, pages 187-204. Interscience Publishers, Inc., New York, N. Y., 1948.
[JR76] R. E. Jamison and W. H. Ruckle. Factoring absolutely convergent series. Math. Ann., 224(2):143-148, 1976.
[JS82a] W. B. Johnson and G. Schechtman. On subspaces of $L_{1}$ with maximal distances to Euclidean space. In Proceedings of research workshop on Banach space theory (Iowa City, Iowa, 1981), pages 83-96. Univ. Iowa, Iowa City, IA, 1982.
[JS82b] I. Joó and L. L. Stachó. Generalization of an inequality of G. Pólya concerning the eigenfrequences of vibrating bodies. Publ. Inst. Math. (Beograd) (N.S.), 31(45):65-72, 1982.
[JS01] W. B. Johnson and G. Schechtman. Finite dimensional subspaces of $L_{p}$. In Handbook of the geometry of Banach spaces, Vol. I, pages 837-870. North-Holland, Amsterdam, 2001. doi:10.1016/S1874-5849(01)80021-8.
[Kah64] J.-P. Kahane. Sur les sommes vectorielles $\sum \pm u_{n}$. C. R. Acad. Sci. Paris, 259:2577-2580, 1964.
[Kal04] N. J. Kalton. Spaces of Lipschitz and Hölder functions and their applications. Collect. Math., 55(2):171-217, 2004.
[Kal08] N. J. Kalton. The complemented subspace problem revisited. Studia Math., 188(3):223-257, 2008.
[Kal12] N. J. Kalton. The uniform structure of Banach spaces. Math. Ann., 354(4):1247-1288, 2012.
[Kaš77] B. S. Kašin. The widths of certain finite-dimensional sets and classes of smooth functions. Izv. Akad. Nauk SSSR Ser. Mat., 41(2):334-351, 478, 1977.
[Kec95] A. S. Kechris. Classical descriptive set theory, volume 156 of Graduate Texts in Mathematics. Springer-Verlag, New York, 1995.
[Kir34] M. D. Kirszbraun. Über die zusammenziehenden und Lipschitzchen Transformationen. Fundam. Math., 22:77-108, 1934.
[KLMN05] R. Krauthgamer, J. R. Lee, M. Mendel, and A. Naor. Measured descent: a new embedding method for finite metrics. Geom. Funct. Anal., 15(4):839-858, 2005.
[KM73] M. I. Kadec and B. S. Mitjagin. Complemented subspaces in Banach spaces. Uspehi Mat. Nauk, 28(6(174)):77-94, 1973.
[KMS98] D. Karger, R. Motwani, and M. Sudan. Approximate graph coloring by semidefinite programming. J. ACM, 45(2):246265, 1998.
[Koz05] M. Kozdoba. Extension of Banach space valued Lipschitz functions. Master's thesis, Technion-Israel Institute of Technology, 2005.
[KP20] Z. Kabluchko and J. Prochno. The maximum entropy principle and volumetric properties of Orlicz balls, 2020. Preprint, available athttps://arxiv.org/abs/2007.05247
[KPR93] P. N. Klein, S. A. Plotkin, and S. Rao. Excluded minors, network decomposition, and multicommodity flow. In S. R. Kosaraju, D. S. Johnson, and A. Aggarwal, editors, Proceedings of the Twenty-Fifth Annual ACM Symposium on Theory of Computing, May 16-18, 1993, San Diego, CA, USA, pages 682-690. ACM, 1993. ISBN 0-89791-591-7.
[KPT20] Z. Kabluchko, J. Prochno, and C. Thäle. Exact asymptotic volume and volume ratio of Schatten unit balls. J. Approx. Theory, 257:105457, 13, 2020.
[Kra26] E. Krahn. Über Minimaleigenschaften der Kugel in drei und mehr Dimensionen. Acta Univ. Dorpat A 9, 1-44 (1926)., 1926.
[KRN65] K. Kuratowski and C. Ryll-Nardzewski. A general theorem on selectors. Bull. Acad. Polon. Sci. Sér. Sci. Math. Astronom. Phys., 13:397-403, 1965.
[KRZ04] A. Koldobsky, D. Ryabogin, and A. Zvavitch. Projections of convex bodies and the Fourier transform. Israel J. Math., 139:361-380, 2004.
[KS89] S. a. Kwapień and C. Schütt. Some combinatorial and probabilistic inequalities and their application to Banach space theory. II. Studia Math., 95(2):141-154, 1989.
[Kwa72] S. Kwapień. Isomorphic characterizations of inner product spaces by orthogonal series with vector valued coefficients. Studia Math., 44:583-595, 1972. Collection of articles honoring the completion by Antoni Zygmund of 50 years of scientific activity, VI.
[Kyr00] I. Kyrezi. On the entropy of the convex hull of finite sets. Proc. Amer. Math. Soc., 128(8):2393-2403, 2000.
[Lew78] D. R. Lewis. Finite dimensional subspaces of $L_{p}$. Studia Math., 63(2):207-212, 1978.
[Lew79] D. R. Lewis. Ellipsoids defined by Banach ideal norms. Mathematika, 26(1):18-29, 1979.
[Lim18] Y. Lim. Contractive barycentric maps and $L^{1}$ ergodic theorems on the cone of positive definite matrices. J. Math. Anal. Appl., 459(1):291-306, 2018.
[Lin64] J. Lindenstrauss. On nonlinear projections in Banach spaces. Michigan Math. J., 11:263-287, 1964.
[LMN05] J. R. Lee, M. Mendel, and A. Naor. Metric structures in $L_{1}$ : dimension, snowflakes, and average distortion. European J. Combin., 26(8):1180-1190, 2005.
[LMSS07] N. Linial, S. Mendelson, G. Schechtman, and A. Shraibman. Complexity measures of sign matrices. Combinatorica, 27(4):439-463, 2007.
[LN03] J. R. Lee and A. Naor. Metric decomposition, smooth measures, and clustering, 2003. Unpublished manuscript, available on request.
[LN04a] J. R. Lee and A. Naor. Absolute Lipschitz extendability. C. R. Math. Acad. Sci. Paris, 338(11):859-862, 2004.
[LN04b] J. R. Lee and A. Naor. Embedding the diamond graph in $L_{p}$ and dimension reduction in $L_{1}$. Geom. Funct. Anal., 14(4):745-747, 2004.
[LN05] J. R. Lee and A. Naor. Extending Lipschitz functions via random metric partitions. Invent. Math., 160(1):59-95, 2005.
[LO99] R. Latała and K. Oleszkiewicz. Gaussian measures of dilatations of convex symmetric sets. Ann. Probab., 27(4):19221938, 1999.
[Loz69] G. J. Lozanovskiĭ. Certain Banach lattices. Sibirsk. Mat. Ž., 10:584-599, 1969.
[LPS00] U. Lang, B. Pavlović, and V. Schroeder. Extensions of Lipschitz maps into Hadamard spaces. Geom. Funct. Anal., 10(6):1527-1553, 2000.
[LR88] F. T. Leighton and S. Rao. An approximate max-flow min-cut theorem for uniform multicommodity flow problems with applications to approximation algorithms. In 29th Annual Symposium on Foundations of Computer Science, White Plains, New York, USA, 24-26 October 1988, pages 422-431. IEEE Computer Society, 1988.
[LS91] N. Linial and M. E. Saks. Decomposing graphs into regions of small diameter. In A. Aggarwal, editor, Proceedings of the Second Annual ACMISIGACT-SIAM Symposium on Discrete Algorithms, 28-30 January 1991, San Francisco, California., pages 320-330. ACM/SIAM, 1991.
[LS98] J. Luukkainen and E. Saksman. Every complete doubling metric space carries a doubling measure. Proc. Amer. Math. Soc., 126(2):531-534, 1998.
[LS05] U. Lang and T. Schlichenmaier. Nagata dimension, quasisymmetric embeddings, and Lipschitz extensions. Int. Math. Res. Not., (58):3625-3655, 2005.
[LT71] J. Lindenstrauss and L. Tzafriri. On the complemented subspaces problem. Israel J. Math., 9:263-269, 1971.
[LT77] J. Lindenstrauss and L. Tzafriri. Classical Banach spaces. I. Springer-Verlag, Berlin-New York, 1977. ISBN 3-540-08072-4. Sequence spaces, Ergebnisse der Mathematik und ihrer Grenzgebiete, Vol. 92.
[LT79] J. Lindenstrauss and L. Tzafriri. Classical Banach spaces. II, volume 97 of Ergebnisse der Mathematik und ihrer Grenzgebiete [Results in Mathematics and Related Areas]. Springer-Verlag, Berlin-New York, 1979. ISBN 3-540-08888-1. Function spaces.
[LT91] M. Ledoux and M. Talagrand. Probability in Banach spaces, volume 23 of Ergebnisse der Mathematik und ihrer Grenzgebiete (3) [Results in Mathematics and Related Areas (3)]. Springer-Verlag, Berlin, 1991. ISBN 3-540-52013-9. doi: 10.1007/978-3-642-20212-4. Isoperimetry and processes.
[Lud02] M. Ludwig. Projection bodies and valuations. Adv. Math., 172(2):158-168, 2002.
[Lud05] M. Ludwig. Minkowski valuations. Trans. Amer. Math. Soc., 357(10):4191-4213, 2005.
[Lus72] N. Lusin. Leçons sur les ensembles analytiques et leurs applications. Chelsea Publishing Co., New York, 1972. Avec une note de W. Sierpiński, Preface de Henri Lebesgue, Réimpression de l'edition de 1930.
[Lut93] E. Lutwak. Selected affine isoperimetric inequalities. In Handbook of convex geometry, Vol. A, B, pages 151-176. North-Holland, Amsterdam, 1993. doi:10.1016/B978-0-444-89596-7.50010-9.
[Luz17] N. N. Luzin. Sur la classification de M. Baire. C. R. Math. Acad. Sci. Paris, 164:91-94, 1917.
[LW49] L. H. Loomis and H. Whitney. An inequality related to the isoperimetric inequality. Bull. Amer. Math. Soc, 55:961-962, 1949.
[LW97] L. Lefton and D. Wei. Numerical approximation of the first eigenpair of the $p$-Laplacian using finite elements and the penalty method. Numer. Funct. Anal. Optim., 18(3-4):389-399, 1997.
[Mah39] K. Mahler. Ein Übertragungsprinzip für konvexe Körper. Časopis Pěst. Mat. Fys., 68:93-102, 1939.
[Mat90] J. Matoušek. Extension of Lipschitz mappings on metric trees. Comment. Math. Univ. Carolin., 31(1):99-104, 1990.
[Mat96] J. Matoušek. On the distortion required for embedding finite metric spaces into normed spaces. Israel J. Math., 93:333-344, 1996.
[Mat02] J. Matoušek. Lectures on discrete geometry, volume 212 of Graduate Texts in Mathematics. Springer-Verlag, New York, 2002. ISBN 0-387-95373-6. doi:10.1007/978-1-4613-0039-7.
[Mau03] B. Maurey. Type, cotype and K-convexity. In Handbook of the geometry of Banach spaces, Vol. 2, pages 1299-1332. North-Holland, Amsterdam, 2003. doi:10.1016/S1874-5849(03)80037-2.
[Maz29] S. Mazur. Une remarque sur l'homéomorphie des champs fonctionels. Studia Math., 1:83-85, 1929.
[McS34] E. J. McShane. Extension of range of functions. Bull. Amer. Math. Soc., 40(12):837-842, 1934.
[Mey00] C. Meyer. Matrix analysis and applied linear algebra. Society for Industrial and Applied Mathematics (SIAM), Philadelphia, PA, 2000. ISBN 0-89871-454-0. doi:10.1137/1.9780898719512. With 1 CD-ROM (Windows, Macintosh and UNIX) and a solutions manual (iv+171 pp.).
[Mil71] V. D. Milman. A new proof of A. Dvoretzky's theorem on cross-sections of convex bodies. Funkcional. Anal. $i$ Priložen., 5(4):28-37, 1971.
[Mil85] V. D. Milman. Almost Euclidean quotient spaces of subspaces of a finite-dimensional normed space. Proc. Amer. Math. Soc., 94(3):445-449, 1985.
[Mil87] V. D. Milman. Some remarks on Urysohn's inequality and volume ratio of cotype 2-spaces. In Geometrical aspects of functional analysis (1985/86), volume 1267 of Lecture Notes in Math., pages 75-81. Springer, Berlin, 1987. doi: 10.1007/BFb0078137.
[Mir60] L. Mirsky. Symmetric gauge functions and unitarily invariant norms. Quart. J. Math. Oxford Ser. (2), 11:50-59, 1960.
[MM96] E. Makai, Jr. and H. Martini. The cross-section body, plane sections of convex bodies and approximation of convex bodies. I. Geom. Dedicata, 63(3):267-296, 1996.
[MM16] K. Makarychev and Y. Makarychev. Metric extension operators, vertex sparsifiers and Lipschitz extendability. Israel J. Math., 212(2):913-959, 2016.
[MN07] M. Mendel and A. Naor. Ramsey partitions and proximity data structures. J. Eur. Math. Soc. (JEMS), 9(2):253-275, 2007.
[MN08] M. Mendel and A. Naor. Metric cotype. Ann. of Math. (2), 168(1):247-298, 2008.
[MN13] M. Mendel and A. Naor. Spectral calculus and Lipschitz extension for barycentric metric spaces. Anal. Geom. Metr. Spaces, 1:163-199, 2013.
[MN15] M. Mendel and A. Naor. Expanders with respect to Hadamard spaces and random graphs. Duke Math. J., 164(8):1471-1548, 2015.
[MN21] M. Mendel and A. Naor. Formulae for absolute extendability operators, 2021. Forthcoming manuscript.
[MNR21] M. Mendel, A. Naor, and Y. Rabani. A weighted Sobolev embedding on the discrete torus and nonexistence of gentle partitions of unity, 2021. Forthcoming manuscript.
[MNT20] J. Matoušek, A. Nikolov, and K. Talwar. Factorization norms and hereditary discrepancy. Int. Math. Res. Not. IMRN, (3):751-780, 2020.
[MP76] B. Maurey and G. Pisier. Séries de variables aléatoires vectorielles indépendantes et propriétés géométriques des espaces de Banach. Studia Math., 58(1):45-90, 1976.
[MP84] M. B. Marcus and G. Pisier. Characterizations of almost surely continuous $p$-stable random Fourier series and strongly stationary processes. Acta Math., 152(3-4):245-301, 1984.
[MP86] V. D. Milman and G. Pisier. Banach spaces with a weak cotype 2 property. Israel J. Math., 54(2):139-158, 1986.
[MP89] V. D. Milman and A. Pajor. Isotropic position and inertia ellipsoids and zonoids of the unit ball of a normed $n$ dimensional space. In Geometric aspects of functional analysis (1987-88), volume 1376 of Lecture Notes in Math., pages 64-104. Springer, Berlin, 1989. doi:10.1007/BFb0090049.
[MPS12] E. Markessinis, G. Paouris, and C. Saroglou. Comparing the $M$-position with some classical positions of convex bodies. Math. Proc. Cambridge Philos. Soc., 152(1):131-152, 2012.
[MS86] V. D. Milman and G. Schechtman. Asymptotic theory of finite-dimensional normed spaces, volume 1200 of Lecture Notes in Mathematics. Springer-Verlag, Berlin, 1986. ISBN 3-540-16769-2. With an appendix by M. Gromov.
[MTJ03] P. Mankiewicz and N. Tomczak-Jaegermann. Quotients of finite-dimensional Banach spaces; random phenomena. In Handbook of the geometry of Banach spaces, Vol. 2, pages 1201-1246. North-Holland, Amsterdam, 2003. doi:10. 1016/S1874-5849(03)80035-9.
[Mül90] D. Müller. A geometric bound for maximal functions associated to convex bodies. Pacific J. Math., 142(2):297-312, 1990.
[MV03] S. Mendelson and R. Vershynin. Entropy and the combinatorial dimension. Invent. Math., 152(1):37-55, 2003.
[MW78] V. D. Mil'man and H. Wolfson. Minkowski spaces with extremal distance from the Euclidean space. Israel J. Math., 29(2-3):113-131, 1978.
[Nao01] A. Naor. A phase transition phenomenon between the isometric and isomorphic extension problems for Hölder functions between $L_{p}$ spaces. Mathematika, 48(1-2):253-271 (2003), 2001.
[Nao07] A. Naor. The surface measure and cone measure on the sphere of $l_{p}^{n}$. Trans. Amer. Math. Soc., 359(3):1045-1079, 2007.
[Nao15] A. Naor. Class notes on Lipschitz extension from finite subsets, 2015. Available at https://web.math.princeton. edu/~naor/homepage\%20files/extension-from-finite.pdf
[Nao17a] A. Naor. Probabilistic clustering of high dimensional norms. In Proceedings of the Twenty-Eighth Annual ACM-SIAM Symposium on Discrete Algorithms, SODA 2017, pages 690-709. 2017. doi:10.1137/1.9781611974782.44.
[Nao17b] A. Naor. A spectral gap precludes low-dimensional embeddings. In 33rd International Symposium on Computational Geometry, volume 77 of LIPIcs. Leibniz Int. Proc. Inform., pages Art. No. 50, 16. Schloss Dagstuhl. Leibniz-Zent. Inform., Wadern, 2017.
[Nao18] A. Naor. Metric dimension reduction: a snapshot of the Ribe program. In Proceedings of the International Congress of Mathematicians-Rio de Janeiro 2018. Vol. I. Plenary lectures, pages 759-837. World Sci. Publ., Hackensack, NJ, 2018.
[Nao2la] A. Naor. An average John theorem. Geom. Topol., 25(4):1631-1717, 2021.
[Nao21b] A. Naor. Impossibility of almost extension. Adv. Math., 384:107761, 2021.
[Nav13] A. Navas. An $L^{1}$ ergodic theorem with values in a non-positively curved space via a canonical barycenter map. Ergodic Theory Dynam. Systems, 33(2):609-623, 2013.
[NPS20] A. Naor, G. Pisier, and G. Schechtman. Impossibility of dimension reduction in the nuclear norm. Discrete Comput. Geom., 63(2):319-345, 2020.
[NR03] I. Newman and Y. Rabinovich. A lower bound on the distortion of embedding planar metrics into Euclidean space. Discrete Comput. Geom., 29(1):77-81, 2003.
[NR17] A. Naor and Y. Rabani. On Lipschitz extension from finite subsets. Israel J. Math., 219(1):115-161, 2017.
[NS07] A. Naor and G. Schechtman. Planar earthmover is not in $L_{1}$. SIAM J. Comput., 37(3):804-826 (electronic), 2007.
[NS16] A. Naor and G. Schechtman. Pythagorean powers of hypercubes. Ann. Inst. Fourier (Grenoble), 66(3):1093-1116, 2016.
[NS21a] A. Naor and G. Schechtman. Lipschitz almost-extension and nonexistence of uniform embeddings of balls in Schatten classes, 2021. Forthcoming manuscript.
[NS21b] A. Naor and G. Schechtman. Obstructions to metric embeddings of Schatten classes, 2021. Forthcoming manuscript.
[NT10] A. Naor and T. Tao. Random martingales and localization of maximal inequalities. J. Funct. Anal., 259(3):731-779, 2010.
[NT12] A. Naor and T. Tao. Scale-oblivious metric fragmentation and the nonlinear Dvoretzky theorem. Israel J. Math., 192(1):489-504, 2012.
[NY21] A. Naor and R. Young. Foliated corona decompositions, 2021. To appear in Acta Math.
[Oht09] S.-i. Ohta. Extending Lipschitz and Hölder maps between metric spaces. Positivity, 13(2):407-425, 2009.
[Ost13] M. I. Ostrovskii. Metric embeddings. Bilipschitz and coarse embeddings into Banach spaces, volume 49 of De Gruyter Studies in Mathematics. De Gruyter, Berlin, 2013.
[Paj83] A. Pajor. Plongement de $l_{1}^{n}$ dans les espaces de Banach complexes. C. R. Acad. Sci. Paris Sér. I Math., 296(17):741-743, 1983.
[Par17] E. Parini. Reverse Cheeger inequality for planar convex sets. J. Convex Anal., 24(1):107-122, 2017.
[Peł68] A. Pełczyński. Linear extensions, linear averagings, and their applications to linear topological classification of spaces of continuous functions. Dissertationes Math. Rozprawy Mat., 58:92, 1968.
[Pet38] B. J. Pettis. Linear functionals and completely additive set functions. Duke Math. J., 4(3):552-565, 1938.
[Pet61] C. M. Petty. Surface area of a convex body under affine transformations. Proc. Amer. Math. Soc., 12:824-828, 1961.
[Pet67] C. M. Petty. Projection bodies. In Proc. Colloquium on Convexity (Copenhagen, 1965), pages 234-241. Kobenhavns Univ. Mat. Inst., Copenhagen, 1967.
[Pet71] C. M. Petty. Isoperimetric problems. pages 26-41, 1971.
[Pis73] G. Pisier. Sur les espaces de Banach qui ne contiennent pas uniformément de $l_{n}^{1}$. C. R. Acad. Sci. Paris Sér. A-B, 277:A991-A994, 1973.
[Pis79] G. Pisier. Sur les espaces de Banach de dimension finie à distance extrémale d'un espace euclidien [d'après V. D. Milman et H. Wolfson]. In Séminaire d'Analyse Fonctionnelle (1978-1979), pages Exp. No. 16, 10. École Polytech., Palaiseau, 1979.
[Pis80a] G. Pisier. Sur les espaces de Banach K-convexes. In Seminar on Functional Analysis, 1979-1980 (French), pages Exp. No. 11, 15. École Polytech., Palaiseau, 1980.
[Pis80b] G. Pisier. Un théorème sur les opérateurs linéaires entre espaces de Banach qui se factorisent par un espace de Hilbert. Ann. Sci. École Norm. Sup. (4), 13(1):23-43, 1980.
[Pis81] G. Pisier. Remarques sur un résultat non publié de B. Maurey. In Seminar on Functional Analysis, 1980-1981, pages Exp. No. V, 13. École Polytech., Palaiseau, 1981.
[Pis82] G. Pisier. Holomorphic semigroups and the geometry of Banach spaces. Ann. of Math. (2), 115(2):375-392, 1982.
[Pis83] G. Pisier. On the dimension of the $l_{p}^{n}$-subspaces of Banach spaces, for $1 \leqslant p<2$. Trans. Amer. Math. Soc., 276(1):201211, 1983.
[Pis86] G. Pisier. Factorization of linear operators and geometry of Banach spaces, volume 60 of CBMS Regional Conference Series in Mathematics. Published for the Conference Board of the Mathematical Sciences, Washington, DC; by the American Mathematical Society, Providence, RI, 1986. ISBN 0-8218-0710-2. doi:10.1090/cbms/060.
[Pis88] G. Pisier. Weak Hilbert spaces. Proc. London Math. Soc. (3), 56(3):547-579, 1988.
[Pis89] G. Pisier. The volume of convex bodies and Banach space geometry, volume 94 of Cambridge Tracts in Mathematics. Cambridge University Press, Cambridge, 1989.
[Pis96] G. Pisier. Dvoretzky's theorem for operator spaces. Houston J. Math., 22(2):399-416, 1996.
[PR98] D. Peleg and E. Reshef. Deterministic polylog approximation for minimum communication spanning trees (extended abstract). In Automata, languages and programming (Aalborg, 1998), volume 1443 of Lecture Notes in Comput. Sci., pages 670-681. Springer, Berlin, 1998.
[PS51] G. Pólya and G. Szegö. Isoperimetric Inequalities in Mathematical Physics. Annals of Mathematics Studies, No. 27. Princeton University Press, Princeton, N. J., 1951.
[PW98] A. Pietsch and J. Wenzel. Orthonormal systems and Banach space geometry, volume 70 of Encyclopedia of Mathematics and its Applications. Cambridge University Press, Cambridge, 1998. ISBN 0-521-62462-2. doi:10.1017/ CBO9780511526145.
[Rad19] H. Rademacher. Über partielle und totale differenzierbarkeit von Funktionen mehrerer Variabeln und über die Transformation der Doppelintegrale. Math. Ann., 79(4):340-359, 1919.
[Rao99] S. Rao. Small distortion and volume preserving embeddings for planar and Euclidean metrics. In Proceedings of the Fifteenth Annual Symposium on Computational Geometry (Miami Beach, FL, 1999), pages 300-306 (electronic). ACM, New York, 1999. doi:10.1145/304893.304983.
[Reg13] O. Regev. Entropy-based bounds on dimension reduction in LL. Israel J. Math., 195(2):825-832, 2013.
[Rei86] S. Reisner. Zonoids with minimal volume-product. Math. Z., 192(3):339-346, 1986.
[Rie27] M. Riesz. Sur les maxima des formes bilinéaires et sur les fonctionnelles linéaires. Acta Math., 49(3-4):465-497, 1927.
[Rog50] C. A. Rogers. A note on coverings and packings. J. London Math. Soc., 25:327-331, 1950.
[RR91a] S. T. Rachev and L. Rüschendorf. Approximate independence of distributions on spheres and their stability properties. Ann. Probab., 19(3):1311-1337, 1991.
[RR91b] M. M. Rao and Z. D. Ren. Theory of Orlicz spaces, volume 146 of Monographs and Textbooks in Pure and Applied Mathematics. Marcel Dekker, Inc., New York, 1991. ISBN 0-8247-8478-2.
[Rut65] D. Rutovitz. Some parameters associated with finite-dimensional Banach spaces. J. London Math. Soc., 40:241-255, 1965.
[RV06] M. Rudelson and R. Vershynin. Combinatorics of random processes and sections of convex bodies. Ann. of Math. (2), 164(2):603-648, 2006.
[Rya02] R. A. Ryan. Introduction to tensor products of Banach spaces. Springer Monographs in Mathematics. Springer-Verlag London, Ltd., London, 2002. ISBN 1-85233-437-1. doi:10.1007/978-1-4471-3903-4.
[San49] L. A. Santaló. An affine invariant for convex bodies of $n$-dimensional space. Portugal. Math., 8:155-161, 1949.
[Sch78] C. Schütt. Unconditionality in tensor products. Israel J. Math., 31(3-4):209-216, 1978.
[Sch82] C. Schütt. On the volume of unit balls in Banach spaces. Compositio Math., 47(3):393-407, 1982.
[Sch89] C. Schütt. The isoperimetric quotient and some classical Banach spaces. Israel J. Math., 67(1):43-61, 1989.
[Sch92] M. Schmuckenschläger. The distribution function of the convolution square of a convex symmetric body in $\mathbf{R}^{n}$. Israel J. Math., 78(2-3):309-334, 1992.
[Sch95] M. Schmuckenschläger. Petty's projection inequality and Santalo's affine isoperimetric inequality. Geom. Dedicata, 57(3):285-295, 1995.
[Sch14] R. Schneider. Convex bodies: the Brunn-Minkowski theory, volume 151 of Encyclopedia of Mathematics and its Applications. Cambridge University Press, Cambridge, expanded edition, 2014. ISBN 978-1-107-60101-7.
[She64] G. C. Shephard. Shadow systems of convex sets. Israel J. Math., 2:229-236, 1964.
[Sil73] E. Silverman. Lower semicontinuity of parametric integrals. Trans. Amer. Math. Soc., 175:499-508, 1973.
[Sim79] B. Simon. Trace ideals and their applications, volume 35 of London Mathematical Society Lecture Note Series. Cambridge University Press, Cambridge-New York, 1979. ISBN 0-521-22286-9.
[Sob41] A. Sobczyk. Projections in Minkowski and Banach spaces. Duke Math. J., 8:78-106, 1941.
[SR84] J. Saint-Raymond. Le volume des idéaux d'opérateurs classiques. Studia Math., 80(1):63-75, 1984.
[Sri98] S. M. Srivastava. A course on Borel sets, volume 180 of Graduate Texts in Mathematics. Springer-Verlag, New York, 1998.
[ST89] S. J. Szarek and M. Talagrand. An "isomorphic" version of the Sauer-Shelah lemma and the Banach-Mazur distance to the cube. In Geometric aspects of functional analysis (1987-88), volume 1376 of Lecture Notes in Math., pages 105-112. Springer, Berlin, 1989. doi:10.1007/BFb0090050.
[STJ80] S. a. Szarek and N. Tomczak-Jaegermann. On nearly Euclidean decomposition for some classes of Banach spaces. Compositio Math., 40(3):367-385, 1980.
[Stu03] K.-T. Sturm. Probability measures on metric spaces of nonpositive curvature. In Heat kernels and analysis on manifolds, graphs, and metric spaces (Paris, 2002), volume 338 of Contemp. Math., pages 357-390. Amer. Math. Soc., Providence, RI, 2003. doi:10.1090/conm/338/06080.
[SW83] R. Schneider and W. Weil. Zonoids and related topics. In Convexity and its applications, pages 296-317. Birkhäuser, Basel, 1983.
[SZ90] G. Schechtman and J. Zinn. On the volume of the intersection of two $L_{p}^{n}$ balls. Proc. Amer. Math. Soc., 110(1):217-224, 1990.
[Sza78] S. a. J. Szarek. On Kashin's almost Euclidean orthogonal decomposition of $l_{n}^{1}$. Bull. Acad. Polon. Sci. Sér. Sci. Math. Astronom. Phys., 26(8):691-694, 1978.
[Sza90] S. a. J. Szarek. Spaces with large distance to $l_{\infty}^{n}$ and random matrices. Amer. J. Math., 112(6):899-942, 1990.
[Sza91] S. J. Szarek. On the geometry of the Banach-Mazur compactum. In Functional analysis (Austin, TX, 1987/1989), volume 1470 of Lecture Notes in Math., pages 48-59. Springer, Berlin, 1991.
[Ta192] M. Talagrand. Type, infratype and the Elton-Pajor theorem. Invent. Math., 107(1):41-59, 1992.
[Ta195] M. Talagrand. Embedding of $l_{k}^{\infty}$ and a theorem of Alon and Milman. In Geometric aspects of functional analysis (Israel, 1992-1994), volume 77 of Oper. Theory Adv. Appl., pages 289-293. Birkhäuser, Basel, 1995.
[Tay47] A. E. Taylor. A geometric theorem and its application to biorthogonal systems. Bull. Amer. Math. Soc., 53:614-616, 1947.
[Tho48] G. O. Thorin. Convexity theorems generalizing those of M. Riesz and Hadamard with some applications. Comm. Sem. Math. Univ. Lund [Medd. Lunds Univ. Mat. Sem.], 9:1-58, 1948.
[Tik19] K. Tikhomirov. On the Banach-Mazur distance to cross-polytope. Adv. Math., 345:598-617, 2019.
[TJ79] N. Tomczak-Jaegermann. Computing 2-summing norm with few vectors. Ark. Mat., 17(2):273-277, 1979.
[TJ89] N. Tomczak-Jaegermann. Banach-Mazur distances and finite-dimensional operator ideals, volume 38 of Pitman Monographs and Surveys in Pure and Applied Mathematics. Longman Scientific \& Technical, Harlow; copublished in the United States with John Wiley \& Sons, Inc., New York, 1989.
[TJ80] N. Tomczak-Jaegermann. Finite-dimensional subspaces of uniformly convex and uniformly smooth Banach lattices and trace classes $S_{p}$. Studia Math., 66(3):261-281, 1979/80.
[Tri49] F. Tricomi. Sulle funzioni di Bellel di ordine e argomento pressochè uguali. Atti Accad. Sci. Torino Cl. Sci. Fis. Mat. Natur., 83:3-20, 1949.
[Vil09] C. Villani. Optimal transport, volume 338 of Grundlehren der Mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences]. Springer-Verlag, Berlin, 2009. ISBN 978-3-540-71049-3. doi:10.1007/ 978-3-540-71050-9. Old and new.
[VK87] A. L. Vol'berg and S. V. Konyagin. On measures with the doubling condition. Izv. Akad. Nauk SSSR Ser. Mat., 51(3):666-675, 1987.
[Wag77] D. H. Wagner. Survey of measurable selection theorems. SIAM J. Control Optimization, 15(5):859-903, 1977.
[Wea99] N. Weaver. Lipschitz algebras. World Scientific Publishing Co., Inc., River Edge, NJ, 1999. ISBN 981-02-3873-8. doi: 10.1142/4100.
[Wey39] H. Weyl. The Classical Groups. Their Invariants and Representations. Princeton University Press, Princeton, N.J., 1939.
[Whi34] H. Whitney. Analytic extensions of differentiable functions defined in closed sets. Trans. Amer. Math. Soc., 36(1):6389, 1934.
[Woj91] P. Wojtaszczyk. Banach spaces for analysts, volume 25 of Cambridge Studies in Advanced Mathematics. Cambridge University Press, Cambridge, 1991.
[WW62] E. T. Whittaker and G. N. Watson. A course of modern analysis. An introduction to the general theory of infinite processes and of analytic functions: with an account of the principal transcendental functions. Cambridge University Press, New York, 1962. Fourth edition. Reprinted.
[Zon02] C. Zong. From deep holes to free planes. Bull. Amer. Math. Soc. (N.S.), 39(4):533-555 (electronic), 2002.
Mathematics Department, Princeton University,, Fine Hall, Washington Road, Princeton, NJ 08544-1000, USA
Email address: naor@math. princeton. edu


[^0]:    ${ }^{1}$ We are suppressing here measurability issues that are addressed in Section 1.7 and Section 3.1
    ${ }^{2}$ We use the following conventions for asymptotic notation, in addition to the usual $O(\cdot), o(\cdot), \Omega(\cdot)$ notation. Given $a, b>0$, by writing $a \lesssim b$ or $b \gtrsim a$ we mean that $a \leqslant C b$ for some universal constant $C>0$, and $a=b$ stands for $(a \lesssim b) \wedge(b \lesssim a)$. If we need to allow for dependence on parameters, we indicate it by subscripts. For example, in the presence of an auxiliary parameter $q$, the notation $a \lesssim q b$ means that $a \leqslant C(q) b$, where $C(q)>0$ may depend only on $q$, and similarly for $a \gtrsim q b$ and $a=q b$.
    ${ }^{3}$ The proof of Theorem 2 artificially gives more such spaces, e.g. $\ell_{\infty}^{n} \oplus \ell_{2}^{n}$, or $\ell_{\infty}^{n} \oplus \mathbf{X}$ for any normed space $\mathbf{X}$ with $\operatorname{dim}(\mathbf{X}) \leqslant \sqrt{n}$.

[^1]:    ${ }^{4}$ Our approach also pertains to subsets of normed spaces, e.g. we will prove that for any $p \in[1, \infty], n \in \mathbb{N}$ and $r \in\{1, \ldots, n\}$, the separation modulus of the set of $n$-by- $n$ matrices of rank at most $r$, equipped with the Schatten-von Neumann- $p$ norm, is equal up to lower order factors to $\max \left\{\sqrt{r}, r^{1 / p}\right\} \sqrt{n}$, which is new even in the Euclidean (Hilbert-Schmidt) setting $p=2$. However, for the purpose of this initial overview we will restrict attention to bounds for the entire space $\mathbf{X}$.
    ${ }^{5}$ By Lud02 Lud05 the mapping that assigns a convex body $K \subseteq \mathbb{R}^{n}$ to its projection body $\Pi K$ is characterized axiomatically as the unique (up to scaling) translation-invariant $S L_{n}(\mathbb{R})$-contravariant Minkowski valuation.

[^2]:    ${ }^{6}$ The assertion that $\mathrm{e}(M)$ is a bi-Lipschitz invariant refers to the fact that the definition immediately implies that if $\left(\eta\right.$, $\left.d_{\eta}\right)$ is another metric space into which $(m, d m)$ admits a bi-Lipschitz embedding, then e $(M) \leqslant \mathrm{c} \eta(M) \mathrm{e}(\eta)$.

[^3]:    ${ }^{7}$ Our proof of the lower bound on $\mathrm{e}(\mathbf{X})$ of Theorem 13 shows that this supremum is at least $\frac{1}{12}$; see equation 140.

[^4]:    ${ }^{8}$ We recommend checking that the analogous stabilization argument does not lead to a similar self-improvement phenomenon in Conjecture 9 Conjecture 10 and Corollary 33 the computations in Section 4 of MPS12 are relevant for this purpose.

[^5]:    ${ }^{9}$ Alternatively, evr $(\mathbf{X}) \lesssim 1$ can be justified by writing $\mathbf{X}=\Pi \mathbf{Z}$ for some normed space $\mathbf{Z}=\left(\mathbb{R}^{n},\|\cdot\|_{\mathbf{Z}}\right)$ (using Bol69 SW83) $)$, and then applying the bound (85) that we derived above (this even demonstrates that the external volume ratio of $\Pi \mathbf{Z}$ is $O(1)$ when $\mathbf{Z}$ is in minimum surface area position rather when $\mathbf{Z}$ is in Löwner position). Actually, the sharp bound $\operatorname{evr}(\mathbf{X}) \leqslant \operatorname{evr}\left(\ell_{\infty}^{n}\right)$ holds, as seen by combining Bal91c Theorem 6] with Reisner's theorem Rei86 that the Mahler conjecture Mah39 holds for zonoids.

[^6]:    ${ }^{10}$ In Naol7a we called the same quantity the "modulus of separated decomposability."

[^7]:    ${ }^{11}$ This assertion is standard, here is a quick sketch. Take a $\delta$-net $n$ of the unit sphere of $\mathbf{X}^{*}$ for a sufficiently small universal
     2-embedding (say, for $d=1 / 10$ ) of $\mathbf{X}$ into $\ell_{\infty}^{n}$ provided $\log n$ is at least a sufficiently large universal constant multiple of dim( $\mathbf{X}$ ).

[^8]:    ${ }^{12}$ By TJ79], if one does not mind losing a universal constant factor in 155, then one could take $m=n$ here, but for the purpose of the ensuing reasoning it suffices to use the much simpler result FLM77, Lemma 6.1].

[^9]:    ${ }^{13}$ Note for later use that if $\mu, v \in \mathrm{M}_{1}^{+}(m)$ satisfy $\mu(m)=v(m)$, then $\mu-v \in \mathrm{M}_{1}^{0}(m)$ and $\|\mu-v\|_{\mathrm{W}_{1}}(m)=\mathrm{W}_{1}(\mu, v)$. For a standard justification of the latter assertion, see e.g. the simple deduction of equation (2.2) in NS07.

[^10]:    ${ }^{14}$ The subset $\mathcal{A}_{n}$ can be taken to be any $\varepsilon_{n}$-net of the unit sphere of $\mathbf{V}_{n}$, for any $\varepsilon_{n} \lesssim n^{-3 / 2}$. Note, however, that the bound that follows from (JL84 (and also MM16. Appendix C]) is $\varepsilon_{n} \lesssim n^{-2}$, and this suffices for the present purposes; see NR17 Theorem 23] for the above stated weaker requirement from $\varepsilon_{n}$.
    ${ }^{15}$ And, it would show that for arbitrarily large $k \in \mathbb{N}$ there exist a metric space ( $m, d m$ ), a Banach space $\left(\mathbf{Z},\|\cdot\|_{\mathbf{Z}}\right)$ and a subset $\mathcal{S} \subseteq m$ with $|\mathcal{S}|=k$ such that $\mathrm{e}_{\mathrm{conv}}(m, \mathcal{S} ; \mathbf{Z}) / \mathrm{e}_{\text {span }}(m, \mathcal{S} ; \mathbf{Z}) \gtrsim \sqrt{(\log k) / \log \log k}$. It would then remain an interesting open question to determine the largest possible asymptotic dependence on $k$ here.

[^11]:    ${ }^{16}$ As $\mathbf{X}_{m}$ is canonically positioned, this holds as an equality, but for the present purposes we just need the stated inclusion.

