# Approximation of zonoids by zonotopes $\left({ }^{1}\right)$ 

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## 1. Introduction

Consider the Euclidean ball $\boldsymbol{B}^{n}$ in $\mathbf{R}^{n}$. It is well known that $\boldsymbol{B}^{n}$ can be approximated in the Hausdorff metric by a sum of segments. Given $\varepsilon>0$, what is the number $N$ needed so that the Hausdorff distance between $B^{n}$ and a sum of segments $\Sigma_{j=1}^{N} I_{j}$ is less than $\varepsilon$ ? It is quite clear that $N=\exp (c(\varepsilon) n)$ for a suitable $c(\varepsilon)$ will suffice. The surprising fact is that actually $N=c(\varepsilon) n$ will do. This was proved in [F.L.M.]. In this paper we show that this fact is not a special property of $B^{n}$ but that essentially the same holds for any convex body in $\mathbf{R}^{n}$ which is a limit of a sum of segments. Questions related to this topic have been studied in the literature till now mainly in the framework of Banach space theory. Also the main tools we use in this paper are taken from Banach space theory. In order to make the paper accessible to experts in convexity theory as well as those in Banach space theory we include in the paper somewhat more than the usual amount of background material.

The introduction is divided into two parts. We first explain the geometric problem and state the main results. We then pass on to the functional analytic formulation, survey the history of the problem and explain the contents of the various sections of this paper.

A zonotope in $\mathbf{R}^{n}$ is a polytope $P$ which is a vector sum of segments $\left\{I_{j}\right\}_{j=1}^{N}$, i.e.,

$$
P=\left\{x ; x=\sum_{j=1}^{N} x_{j}, x_{j} \in I_{j}, j=1, \ldots, N\right\} .
$$

[^0](By a segment we mean a compact one dimensional convex set.) A zonotope always has a center of symmetry, namely the sum of the centers of the segments $I_{j}$. It is easy to check that every face of a zonotope is again a zonotope and hence all faces of a zonotope have a center of symmetry. Conversely, if $P$ is a polytope all of whose faces have a center of symmetry then $P$ is a zonotope (it actually suffices to check just its two dimensional faces; for a proof of these facts and other facts mentioned below we refer to the expository paper [Bo]). Thus, it is easy to check which polytopes are zonotopes; in $\mathbf{R}^{2}$ all centrally symmetric polytopes (i.e., convex polygons) are zonotopes, in $\mathbf{R}^{3}$ the cube for example is a zonotope while the octahedron is not. A zonoid in $\mathbf{R}^{n}$ is defined to be a limit in the Hausdorff metric of a sequence of zonotopes, and thus in particular is a compact convex set with a center of symmetry. From what was said above, it follows that any compact convex symmetric set in $\mathbf{R}^{2}$ is a zonoid. For $n \geqslant 3$ the unit ball in $l_{p}^{n}$ i.e.,
$$
\left\{x=\left(t_{1}, t_{2}, \ldots, t_{n}\right) ;\|x\|=\left(\sum_{i=1}^{n}\left|t_{i}\right|^{p}\right)^{1 / p} \leqslant 1\right\}
$$
is a zonoid if and only if $p \geqslant 2$. In particular, the usual Euclidean unit ball $B^{n}$ in $l_{2}^{n}$ is a zonoid. A polytope is a zonoid if and only if it is already a zonotope. The zonoids arise naturally in many contexts. For example the zonoids are exactly the ranges of atom free vector (i.e., $\mathbf{R}^{n}$ valued) measures.

Clearly, an affine image of a zonotope (zonoid) is itself a zonotope (resp. zonoid). In general the polar body of a zonoid fails to be a zonoid. It follows from a result of Grothendieck that there is an absolute constant $\gamma$ so that if $B$ is a convex body in $\mathbf{R}^{n}$ with the origin as its center of symmetry so that both it and its polar

$$
B^{0}=\{y ;|\langle x, y\rangle| \leqslant 1 \text { all } x \in B\}
$$

are zonoids then there is an affine automorphism $T$ of $\mathbf{R}^{n}$ so that

$$
B \subset T\left(B^{n}\right) \subset \gamma B
$$

There are however, for every $n$, convex bodies $B$ in $\mathbf{R}^{n}$ which are not affine images of balls so that both $B$ and its polar are zonoids (cf. [Schn.] and also the recent survey paper [Schn.W]). The Euclidean ball $B^{n}$ plays a central role in the theory of the zonoids; the result of Grothendieck mentioned above is just one indication of this fact.

The problem we study in the present paper is the following: Given a zonoid $B$ in $\mathbf{R}^{n}$
and given $\varepsilon>0$ what is the minimal $N=N(B, \varepsilon)$ so that there is a sum $P_{N}$ of $N$ segments so that

$$
\begin{equation*}
B \subset P_{N} \subset(1+\varepsilon) B \tag{1.1}
\end{equation*}
$$

The $\varepsilon$ measures the distance between $B$ and $P_{N}$. It is related to the Hausdorff distance $\varrho$ between $B$ and $P_{N}$ by the obvious relation

$$
\begin{equation*}
\varepsilon r<\varrho<\varepsilon R \tag{1.2}
\end{equation*}
$$

where $r$ (respectively $R$ ) is the radius of the largest Euclidean ball contained (resp. the smallest ball containing) $B$. This $\varepsilon$ is more natural in our context than $\varrho$ since it makes $N(B, \varepsilon)$ an affine invariant of $B$.

As mentioned already above for $B=B^{n}$ the Euclidean ball, it was proved in [F.L.M.] that (if $\varepsilon<1 / 2$ say)

$$
\begin{equation*}
N\left(B^{n}, \varepsilon\right) \leqslant c \varepsilon^{-2} \log \varepsilon^{-1} \cdot n \tag{1.3}
\end{equation*}
$$

where $c$ is an absolute constant. This estimate was slightly improved in [G] where it is shown that

$$
\begin{equation*}
N\left(B^{n}, \varepsilon\right) \leqslant c \varepsilon^{-2} n \tag{1.4}
\end{equation*}
$$

The main interest in (1.3) (or (1.4)) is in the fact that $N\left(B^{n}, \varepsilon\right)$ depends linearly on $n$. This dependence on $n$ is certainly optimal. As far as the dependence on $\varepsilon$ is concerned it is obviously not optimal for $n=2$ where one easily sees that $N\left(B^{2}, \varepsilon\right) \approx \varepsilon^{-1 / 2}$ as $\varepsilon \rightarrow 0$. It was pointed out in [Be.Mc] that also for $n=3,4$ (1.4) can be improved. It was proved recently that also for $n \geqslant 5$ the dependence on $\varepsilon$ in (1.4) can be improved (cf. [Li]). (') We shall prove however below (in section 6) that for a given $\tau>0$ we have for all sufficiently large $n$ that

$$
N\left(B^{n}, \varepsilon\right) \geqslant c(n) \varepsilon^{-2+r}
$$

for a suitable positive constant $c(n)$ depending on $n$.
The main results in this paper show that the estimate obtained previously for $B^{n}$ and which, as we just remaked, is quite sharp in this case, is valid in a slightly weaker form for an arbitrary zonoid.

[^1]Theorem 1.1. Let B be a zonoid in $\mathbf{R}^{n}$ which is the unit ball of a uniformly convex norm in $\mathbf{R}^{n}$ and let $\tau>0$. Then we have

$$
\begin{equation*}
N(B, \varepsilon) \leqslant C \varepsilon^{-(2+\tau)} n \tag{1.5}
\end{equation*}
$$

where the constant $C=C(\tau, \delta)$ depends only on $r$ and the degree of uniform convexity, i.e., on the constant

$$
\begin{equation*}
\delta=\min \{2-\|x+y\| ;\|x\|,\|y\| \leqslant 1,\|x-y\| \geqslant 1\} \tag{1.6}
\end{equation*}
$$

Where || || is the norm in $\mathbf{R}^{n}$ whose unit ball is $B$.

Again, the main point here is the linear dependence of $N(B, \varepsilon)$ on $n$. For general zonoids in $\mathbf{R}^{n}$ (i.e., without the additional uniform convexity assumption) we have the following result.

Theorem 1.2. For every $\tau>0$ there is a constant $c(\tau)$ so that for any zonoid $B$ in $\mathbf{R}^{n}, n \geqslant 2$,

$$
\begin{equation*}
N(B, \varepsilon) \leqslant c(\tau) \varepsilon^{-(2+\tau)}(\log n)^{3} n . \tag{1.7}
\end{equation*}
$$

The proofs of both theorems give the following additional information on the approximating zonotopes. If $B$ is a zonoid and $P$ a zonotope $\Sigma_{j=1}^{N} I_{j}$ which approximates $B$ up to $\varepsilon$ then there exist $\left\{j_{i}\right\}_{i=1}^{\mathcal{N}}$ where $\tilde{N}$ does not exceed the right hand side of (1.5) (resp. (1.7)) and scalars $\left\{\lambda_{i}\right\}_{i=1}^{N}$ so that $\sum_{i=1}^{N} \lambda_{i} I_{j_{i}}$ approximates $P$ up to $2 \varepsilon$.

Before leaving the geometric part of the introduction, we would like to point out another interpretation of our problem and results. Let $B$ be a convex compact subset of $\mathbf{R}^{n}$ having the origin as an interior point and center of symmetry. For a unit vector $x$ in $\mathbf{R}^{n}$ (taken in the usual Euclidean norm) we denote by $h(x, B)$ the ( $n-1$ )-dimensional volume of the projection of $B$ on the hyperplane orthogonal to $x$. It is well known (cf. [Bon.F.]) that

$$
h(x, K)=\frac{1}{2} \int_{S^{n-1}}|\langle x, u\rangle| d \sigma_{B}(u)
$$

where $S^{n-1}=\partial\left(B^{n}\right)$ and $\sigma_{B}(u)$ is a certain positive measure on $S^{n-1}$ which is derived from the surface measure of $B$. Hence $h(x, K)$ is the restiction to $S^{n-1}$ of a norm on $\mathbf{R}^{n}$. The unit ball of the dual of this norm is a zonoid $Z(B)$ (cf. (1.8) below) which is called the projection body of $B$. The correspondence $B \leftrightarrow Z(B)$ turns out to be a one to one map
from the set of all convex bodies onto the set of all zonoids (we consider only those sets in $\mathbf{R}^{n}$ which have the origin as an interior point and center of symmetry). This classic fact is due to Minkowski and A. D. Aleksandrov (see [Bo] and [Schn.W.] as well as [Bo.F.] for details). It is easy to see and well known that $B$ is a polytope if and only if $Z(B)$ is a zonotope; the number of segments in the representation of $Z(B)$ as a zonotope is equal to half the number of $(n-1)$-dimensional faces of $B$. Thus the problem we consider here can be phrased as follows. Given a centered symmetric convex body $B$ in $\mathbf{R}^{n}$ and $\varepsilon>0$, find the smallest possible $N$ so that there is a polytope $P$ with $N$ faces in $\mathbf{R}^{n}$ for which

$$
Z(P) \subset Z(B) \subset(1+\varepsilon) Z(P)
$$

A recent study of projection bodies which is closely related to the present paper is [B.L.].

We pass now to functional analytic terminology. Note first that by definition the $\operatorname{sum} P$ of $N$ segments is an affine image of the unit cube in $\mathbf{R}^{N}$. Hence, if the origin is the center and an interior point of $P$, the norm induced by $P$ on $\mathbf{R}^{n}$ is a quotient norm of $l_{\infty}^{N}$ (whose unit ball is the $N$-dimensional cube). Consequently the polar $P^{0}$ of $P$ is the unit ball of an $n$ dimensional subspace of $l_{1}^{N}=\left(l_{\infty}^{N}\right)^{*}$. Similarly a convex body $B$ in $\mathbf{R}^{n}$ with center at the origin is a zonoid if and only if the norm $\left\|\|_{*}\right.$ induced by $B^{0}$ is given by

$$
\begin{equation*}
\|x\|_{*}=\int_{S^{n-1}}|\langle x, u\rangle| d \sigma(u) \tag{1.8}
\end{equation*}
$$

where $\sigma(u)$ is a positive measure on the sphere $S^{n-1}$ symmetric with respect to 0 . (The zonoid is a zonotope if and only if the measure $\sigma$ is purely atomic with a finite number of atoms.) In other words $B$ is a zonoid if and only if $B^{0}$ is the unit ball of an $n$ dimensional subspace of $L_{1}(\Omega, \mu)$ for some measure space ( $\Omega, \mu$ ) (as mentioned above we may always take $\Omega=S^{n-1}$, another possible universal choice is $\Omega=[0,1]$ and $\mu=$ the Lebesgue measure; all these elementary facts are explained in detail in [Bo]).

Let us recall the notion of the Banach Mazur distance between two finite dimensional spaces $X$ and $Y$ (of the same dimension)

$$
\begin{equation*}
d(X, Y)=\inf \left\{\|T\|\left\|T^{-1}\right\| ; T \text { linear isomorphism from } X \text { to } Y\right\} \tag{1.9}
\end{equation*}
$$

It follows that $d(X, Y) \leqslant d$ if and only if there is an affine map $T$ so that $B_{X} \subset T B_{Y} \leqslant d B_{X}$ where $B_{X}$ (resp. $B_{Y}$ ) denotes the unit ball of $X$ (resp. $Y$ ). We can now restate our problem as follows.
$\left.{ }^{( }{ }^{*}\right)$ Given an $n$ dimensional subspace $X$ of $L_{1}(0,1)$ and $\varepsilon>0$, what is the smallest $N=N(X, \varepsilon)$ such that there is a subspace $Y$ of $f_{1}^{N}$ with $d(X, Y) \leqslant 1+\varepsilon$.

As mentioned above this question was considered in [F.L.M] and [G] for $X=l_{2}^{n}$. For $X=l_{p}^{n}, 1 \leqslant p<2$ it was proved by Johnson and Schechtman [J.Sche.1] that there is a function $c(\varepsilon, p)$ so that $N\left(l_{p}^{n}, \varepsilon\right) \leqslant c(p, \varepsilon) n$. The dependence of $c(p, \varepsilon)$ on $\varepsilon$ was of the order $\varepsilon^{-p^{\prime}}$ where $p^{\prime}$ is the conjugate exponent to $p$ and thus the result is weaker than what can be deduced from (1.5) in this special case. The proof of [J.Sche.1] depends on a delicate analysis of truncations of $p$-stable variables. The result of [J.Sche.1] was generalized by Pisier [P.1] who considered economical embeddings of $l_{p}^{n}$ spaces into general spaces with a known "stable type $p$ " constant. Schechtman continued this line of research and obtained, using probability distributions related to stable variables some results on general subspaces of $L_{1}(0.1)$ (cf. [Sche.1] or [Sche.2]). In [Sche.3] Schechtman attacked the problem by a new and much simpler method using what is called in probability theory "the empirical distribution method". Using this method he proved that for every subspace $X$ of $L_{1}(0,1)$ of dimension $n$ and every $0<\varepsilon<1 / 2$

$$
\begin{equation*}
N(X, \varepsilon) \leqslant c \varepsilon^{-2} \log \varepsilon^{-1} \cdot n^{2} \tag{1.10}
\end{equation*}
$$

for some absolute constant $c$.
In section 2 below we recall the method of Schechtman [Sche.3]. We do this since our main results are obtained by analyzing and refining the same procedure.

In section 3 we prove a weak version of Theorem 1.1 above in which the dependence of $N(\boldsymbol{B}, \varepsilon)$ in $n$ is already linear but the dependence on $\varepsilon$ is worse than that of (1.7). For $X=l_{p}^{n}$ the dependence on $\varepsilon$ we get in section 3 is essentially the same as that in [J.Sche.1]. The main ingredient we use in section 3 besides the empirical distribution method is a factorization theorem of Pisier [P.2] (or equivalently earlier results of this nature due to Nikishin and Maurey) which shows that we may change the position of a uniformly convex subspace $X$ of $L_{1}(0,1)$ so that in the new position $X$ actually sits "nicely" in $L_{p}(\mathbf{0}, 1)$ for some $p>1$.

In section 4 we prove an estimate on the entropy of the unit ball of an $n$ dimensional subspace of $L_{1}(\mu)$ for some probability measure $\mu$ in the $\left\|\|_{\infty}\right.$ norm. In other words, we compute for a given $r>1$ and for a suitable "good positioning" of an $n$ dimensional subspace $X$ of $L_{1}(\mu)$ the number of balls of radius $r$ in the norm of $L_{\infty}(\mu)$ (and also of $L_{p}(\mu)$ for large finite $p$ ) needed in order to cover the unit ball of $X$ (in the norm induced from $\left.L_{1}(\mu)\right)$. The proof involves Sudakov's inequality and its dual as well as the notion of absolutely summing operators.

In section 5 the methods of section 3 and 4 are combined and used to prove Theorems 1.1 and 1.2.

Section 6 is devoted to the approximation of the usual Euclidean ball $B^{n}$. This section can be read directly after section 2 . In the first part of this section we show that for every compact convex set $K$ in $\mathbf{R}^{n}$ we can approximate up to $\varepsilon$ the ball $B^{n}$ by a Minkowski sum of $N$ rotations of $\lambda K$ where $\lambda$ is a suitable scalar and $N$ is given by the right hand side of (1.3). The second part of section 6 is devoted to the examination of the dependence on $\varepsilon$ of $N\left(B^{n}, \varepsilon\right)$. The main result here (which is obtained by using spherical harmonics) is the fact that $N\left(B^{n}, \varepsilon\right) \geqslant c(n) \varepsilon^{-2(n-1) /(n+2)}$ as $\varepsilon \rightarrow 0$ for fixed $n$. Section 6 concludes that part of the paper which is devoted directly to the geometric problem of approximating zonoids.

All the papers we quoted above in connection with problem (*) deal also with the problem of embedding $n$-dimensional subspaces $X$ of $L_{r}(0,1)$ with $1<r$ into $l_{r}^{N}$ for a suitable small $N$. The answer to this question in the special case $X=l_{2}^{n}$ is given in [F.L.M.]. In section 7 it is shown that, in analogy to the situation for $r=1$, the estimates obtained in [F.L.M.] for $l_{2}^{n}$ are valid in a slightly weaker form for an arbitrary $n$-dimensional subspace of $L_{r}(0,1)$. These results strengthen the previous results in this direction (in [J.Sche.1], [P.1], [Sche.1] and [Sche.3]).

If $X$ is an $n$-dimensional subspace of $L_{r}(0,1)$ which is nicely complemented we may wish to find $Y$ in $l_{r}^{N}$ with say $d(X, Y) \leqslant 2$ so that not only $N$ is small but that $Y$ is also nicely complemented in $l_{r}^{N}$. In section 8 we show, using the empirical distribution method, that this can be done for a general $X$ and the estimate for $N$ obtained is slightly weaker than the estimates obtained in [F.L.M.] for the special case of $X=l_{2}^{n}$. We also apply these results to the siuation studied in [B.Tz.1]. In particular we deduce a result on the Banach space structure of some spaces which arise in harmonic analysis.

The last section of the paper is again concerned with entropy estimates. We obtain a sharper form of estimates proved in section 4. These sharper estimates, in view of their connection with classical questions, are of some independent interest (though they do not yield better estimates for $N(X, \varepsilon))$. These estimates also have applications to questions in approximation theory related to Sobolev inequalities, especially the socalled Bernstein width (cf. [B.G.]).

The main results of this paper were announced in [B.L.M.1].
In this paper we consider only Banach spaces of dimension strictly larger than 1 over the real field. The results and their proofs are valid (with some minor modifications) also in the complex case. There are several universal constants which enter into the estimates below. These constants are denoted by letters like $\gamma, c, C, c_{1}, c_{2}, \ldots$ We
did not distinguish carefully between the different constants, neither did we try to get good estimates for them (the proofs however, yield naturally some estimates). The same letter will be used to denote different universal constants in different parts of the paper. The notation such as $c(\tau, \delta)$ will mean a function depending only on the parameters $\tau$ and $\delta$. The cardinality of a (finite) set $A$ is denoted by $\overline{\bar{A}}$. The characteristic function of a subset $A$ of $\Omega$ is denoted by $\chi_{A}$ (provided $\Omega$ is clear from the context).

## 2. The empirical distribution method

In this section we present the method and main result of Schechtman [Sche.3].
Let $X$ be an $n$-dimensional subspace of $L_{1}(\Omega, \mu)$ where $(\Omega, \mu)$ is some probability space. Let $t_{1}, t_{2}, \ldots, t_{N}$ be $N$ randomly (and independently) chosen points in this measure space. Consider the linear map $T: L_{1}(\mu) \rightarrow l_{1}^{N}$ defined by

$$
\begin{equation*}
T f=N^{-1}\left(f\left(t_{1}\right), f\left(t_{2}\right), \ldots, f\left(t_{N}\right)\right) \tag{2.1}
\end{equation*}
$$

The basic idea of the method is that for $N$ sufficiently large and some (usually most) choices of $\left(t_{1}, t_{2}, \ldots, t_{N}\right)$ the map $T$ is almost an isometry on $X$, i.e.,

$$
\begin{equation*}
\sup \left\{\left|1-N^{-1} \sum_{j=1}^{N}\right| f\left(t_{j}\right)|\mid ; f \in X,\|f\|=1\}\right. \tag{2.2}
\end{equation*}
$$

is small. The tool for verifying that this happens is the following standard classical estimate from probability theory, often called Bernstein's inequality.

Lemma 2.1. Let $\left\{\dot{g}_{j}\right\}_{j=1}^{N}$ be independent random variables with mean 0 on some probability space $(\Omega, \mu)$ which satisfy

$$
\begin{equation*}
\left\|g_{j}\right\|_{1} \leqslant 2, \quad\left\|g_{j}\right\|_{\infty} \leqslant M, \quad 1 \leqslant j \leqslant N \tag{2.3}
\end{equation*}
$$

for some constant $M$. Then for $0<\varepsilon<1$

$$
\begin{equation*}
\operatorname{Prob}\left\{\left|\sum_{j=1}^{N} g_{j}\right| \geqslant \varepsilon N\right\}<2 \exp \left(-\varepsilon^{2} N / 8 M\right) \tag{2.4}
\end{equation*}
$$

Proof. For $-\infty<x \leqslant 1$ we have $e^{x} \leqslant 1+x+x^{2}$. Hence, since $\int g_{j} d \mu=0$ and $\left|g_{j}(t)\right| \leqslant M$ for every $j$ and every $t \in \Omega$ we have if $0<\lambda M \leqslant 1$

$$
\int \exp \left(\lambda g_{j}\right) d \mu \leqslant 1+\lambda^{2} \int g_{j}^{2} d \mu \leqslant 1+\lambda^{2}\left\|g_{j}\right\|_{1}\left\|g_{j}\right\|_{\infty}<\exp \left(2 \lambda^{2} M\right)
$$

By the independence of $\left\{g_{j}\right\}_{j=1}^{N}$ we deduce that

$$
\int \exp \left(\lambda \sum_{j=1}^{N} g_{j}\right) d \mu<\exp \left(2 \lambda^{2} N M\right)
$$

and consequently

$$
\exp (\lambda N \varepsilon) \operatorname{Prob}\left\{\sum_{j=1}^{N} g_{j} \geqslant \varepsilon N\right\}<\exp \left(2 \lambda^{2} N M\right)
$$

Choosing $\lambda=\varepsilon / 4 M$ we certainly fulfil $0<\lambda M<1$ and get

$$
\operatorname{Prob}\left\{\sum_{j=1}^{N} g_{j} \geqslant \varepsilon N\right\}<\exp \left(-\varepsilon^{2} N / 8 M\right)
$$

and this implies (2.4).

Corollary 2.2. Let $(\Omega, \mu)$ be a probability space, let $\mathscr{F}$ be a finite set in $L_{1}(\Omega, \mu)$ so that

$$
\begin{equation*}
\|f\|_{1} \leqslant 1, \quad\|f\|_{\infty} \leqslant M, \quad f \in \mathscr{F} \tag{2.5}
\end{equation*}
$$

for some constant $M$. Let $0<\varepsilon<1$ and let $N$ be an integer so that

$$
\begin{equation*}
2 \overline{\mathscr{F}} \leqslant \exp \left(\varepsilon^{2} N / 8 M\right) \tag{2.6}
\end{equation*}
$$

then there exist $\left\{t_{j}\right\}_{j=1}^{N}$ in $\Omega$ so that

$$
\begin{equation*}
\left|\|f\|_{1}-N^{-1} \sum_{j=1}^{N}\right| f\left(t_{j}\right)|\mid \leqslant \varepsilon, \quad f \in \mathscr{F} . \tag{2.7}
\end{equation*}
$$

Proof. For every $f \in \mathscr{F}$ consider the random variables $\left\{g_{j}\right\}_{j=1}^{N}$ on $\left(\Omega^{N}, \mu^{N}\right)$ defined by

$$
g_{j}(t)=\left|f\left(t_{j}\right)\right|-\int|f| d \mu, \quad j=1, \ldots, N ; \quad t=\left(t_{1}, t_{2}, \ldots, t_{N}\right) \in \Omega^{N}
$$

The $\left\{g_{j}\right\}_{j=1}^{N}$ satisfy the assumptions in Lemma 2.1 and hence by (2.4)

$$
\operatorname{Prob}\left\{t \in \Omega^{N},\left|\|f\|_{1}-N^{-1} \sum_{j=1}^{N}\right| f\left(t_{j}\right)| |>\varepsilon\right\}<2 \exp \left(-\varepsilon^{2} N / 8 M\right)
$$

In view of (2.6) we deduce that there is a $t \in \Omega^{N}$ so that (2.7) holds.

For applying Lemma 2.1 and Corollary 2.2, a good estimate on $\left\|\|_{\infty}\right.$ is needed for the subsets of $L_{1}$ we consider. Such an estimate is provided by the next lemma.

Lemma 2.3. Let $X$ be an n-dimensional subspace of $L_{1}(\Omega, \mu)$ for some probability space $(\Omega, \mu)$. Then there is a linear operator $U: L_{1}(\Omega, \mu) \rightarrow L_{1}(\Omega, v)$ for some probability measure $v$ on $\Omega$ so that $U$ is an isometry on $X$ and

$$
\begin{equation*}
\|f\|_{L_{\infty}(\Omega, v)} \leqslant n\|f\|_{L_{1}(\Omega, v)} \text { for every } f \in U X \tag{2.8}
\end{equation*}
$$

Proof. Let $\left\{f_{i}\right\}_{i=1}^{n}$ be an Auerbach basis in $X$, i.e., an algebraic basis of $X$ so that $\left\|f_{i}\right\|=1$ for all $i$ and so that whenever $f=\sum_{i=1}^{n} a_{i} f_{i}$ is in the unit ball of $X$ then $\left|a_{i}\right| \leqslant 1$ for all $i$ (the existence of such a basis is easy to prove cf. e.g., [L.Tz.1. p.16]). Let $F(t)=$ $\sum_{i=1}^{n}\left|f_{i}(t)\right|$. Then $\|F\|_{1}=n ;|f(t)| \leqslant F(t)$ for every $t \in \Omega$ and every $f$ in the unit ball of $X$. Let $v$ be defined by $d \nu=n^{-1} F d \mu$ and $U: L_{1}(\mu) \rightarrow L_{1}(v)$ by

$$
U g(t)= \begin{cases}n g(t) / F(t) & \text { if } F(t) \neq 0 \\ 0 & \text { if } F(t)=0\end{cases}
$$

With these definitions all the assertions are evident.

The preceding lemma asserts that for an $n$-dimensional subspace $X$ of $L_{1}(\mu)$ we can embed $X$ in $L_{1}(v)$ for another measure so that in the new position we have $\|f\|_{\infty} \leqslant n$ for every $f$ in the unit ball of $X$. The operation of passing from $\mu$ to $v$ is called a change of density operation. If $F(t) \neq 0$ for all $t$ then the map $U$ is an isometry from $L_{1}(\Omega, \mu)$ onto $L_{1}(\Omega, v)$.

A subset $\mathscr{F}$ of a set $D$ is called an $\varepsilon$-net in $D$ for some $\varepsilon>0$ if for every $x \in D$ there in a $y \in \mathscr{F}$ with $\|x-y\| \leqslant \varepsilon$. The following two lemmas are completely standard facts.

Lemma 2.4. The unit ball B of an n-dimensional space (and any subset of it) has for every $\varepsilon>0$ an $\varepsilon$ net of cardinality $\leqslant(1+2 / \varepsilon)^{n}$.

Proof. Consider a maximal subset $\mathscr{F}$ of $B$ with the property that $\|u-v\| \geqslant \varepsilon$ for every $u, v \in \mathscr{F}, u \neq v$. The open balls with centers in $\mathscr{F}$ and radius $\varepsilon / 2$ are mutually disjoint and all are contained in $(1+\varepsilon / 2) B$. By considering the volumes of these balls we get that $\stackrel{\bar{F}}{\mathscr{F}} \leqslant(1+2 / \varepsilon)^{n}$.

Lemma 2.5. Let $T$ be a bounded linear map from a Banach space $X$ into a Banach space $Y$. Let $0<\varepsilon<1$ and assume that for an $\varepsilon$ net $\mathscr{F}$ of the unit ball of $X$

$$
\mid\|T x\|-\|x\| \leqslant \varepsilon, \quad x \in \mathscr{F} .
$$

Then

$$
\begin{equation*}
\mid\|T x\|-\|x\|\|\leqslant \Delta\| x \| \text { for every } x \in X \text { where } \quad \Delta=3 \varepsilon /(1-\varepsilon) \tag{2.9}
\end{equation*}
$$

Proof. Put $\Delta=\sup \{|\|T x\|-\|x\||,\|x\| \leqslant 1\}$. Let $x$ be any point in the unit ball of $X$ and let $y \in \mathscr{F}$ satisfy $\|x-y\| \leqslant \varepsilon$. By the triangle inequality

$$
|\|T x\|-\|x\|| \leqslant\|T y\|-\|y\|+\mid\|T(x-y)\|-\|x-y\|\|+2\| x-y \| .
$$

Hence $\Delta \leqslant \varepsilon+\varepsilon \Delta+2 \varepsilon$ and (2.9) follows.

We now have all the tools needed to prove

Theorem 2.6 (Schechtman [Sche.3]). Let $X$ be an n-dimensional subspace of $L_{1}[0,1]$ and let $0<\varepsilon<1 / 2$. Then there is a subspace $Y$ of $l_{1}^{N}$ with $d(X, Y) \leqslant 1+\varepsilon$ provided that

$$
\begin{equation*}
N \geqslant c \varepsilon^{-2} \log \frac{1}{\varepsilon} \cdot n^{2} \tag{2.10}
\end{equation*}
$$

for some absolute constant c.
Proof. By Lemma 2.3 we may assume that $\|x\|_{\infty} \leqslant n$ for every $x$ in the unit ball of $X$. Let $\mathscr{F}$ be an $\varepsilon$-net in the unit ball of $X$ of cardinality $\leqslant(3 / \varepsilon)^{n}$ (Lemma 2.4). By Corollary 2.2 if

$$
\log 2+n \log (3 / \varepsilon) \leqslant \varepsilon^{2} N / 8 n
$$

i.e., if (2.10) holds for a suitable $c$, the map $T$ from $X$ into $l_{1}^{N}$ defined by (2.1) satisfies $|\|T x\|-\|x\||<\varepsilon$ for every $x \in \mathscr{F}$. Hence by Lemma 2.5,

$$
d(X, T X) \leqslant(1+\Delta)(1-\Delta)^{-1} \leqslant 1+8 \varepsilon
$$

if $\varepsilon$ is small enough.

Our goal in the coming three sections is to improve the estimate (2.10), i.e., to pass from $n^{2}$ to $n$ (or almost $n$ ). The lemmas used above by themselves are sharp, i.e.,
cannot be improved. For example in Lemma 2.3 we cannot replace $n$ by any smaller number. If $X=l_{1}^{n}$ then in any isometric embedding of $X$ into $L_{1}(\Omega, v)$ the unit vectors of $X$ have to map onto disjointly supported functions and thus since $v$ is a probability measure at least one of the unit vectors of $X$ will have an image whose $\left\|\|_{\infty}\right.$ norm is at least $n$. The point which makes the improvement possible, however, is the exploitation of the fact that the set $\mathscr{F}$ to which we applied Corollary 2.2 is not an arbitrary set of cardinality $(3 / \varepsilon)^{n}$ in $L_{1}$ but a set sitting in a low dimensional subspace. Thus while we cannot avoid the situation that some members of $\mathscr{F}$ have an $\left\|\|_{\infty}\right.$ norm equal to $n$, we can manage the data so that this happens only to relatively few elements or alternatively that a large value is attained only for a very small subset of $\Omega$ (substantially smaller than what follows just from the condition that $\|f\|_{1}=1$ ).

## 3. The uniformly convex case-the iteration procedure

In this section we shall treat uniformly convex spaces $X$ and prove a weak version of Theorem 1.1. For the purpose of proving this result we have to measure quantitatively the degree of uniform convexity. In the introduction we defined the constant $\delta$ in (1.6). This $\delta$ is perhaps the easiest measure of uniform convexity to understand geometrically. However, for our purposes (and many other contexts as well) the notion of "type" is more useful.

The Banach space $X$ is said to be of type $1 \leqslant p \leqslant 2$ if there is a constant $\alpha$ so that for every choice of $\left\{x_{i}\right\}_{i=1}^{m}, m=1,2, \ldots$ in $X$

$$
\begin{equation*}
\int_{0}^{1}\left\|\sum_{i=1}^{m} r_{i}(t) x_{i}\right\| d t \leqslant \alpha\left(\sum_{i=1}^{m}\left\|x_{i}\right\|^{p}\right)^{1 / p}, \tag{3.1}
\end{equation*}
$$

where $\left\{r_{i}\right\}_{i=1}^{\infty}$ denote the Rademacher functions. The quantity on the left hand side of (3.1) is just the average of $\left\|\sum_{i=1}^{m} \pm x_{i}\right\|$ over all the possible $2^{m}$ choices of signs. The smallest $\alpha$ for which (3.1) holds is called the type $p$ constant of $X$ and is denoted by $T_{p}(X)$. Evidently $T_{1}(X)=1$ for every $X$. Every finite dimensional $X$ is type $p$ for every $1 \leqslant p \leqslant 2$ (this, as well as the fact that $p$ can never be taken $>2$, follows immediately from the generalized parallelogram identity in Hilbert space, i.e., the identity

$$
\int\left\|\sum_{i=1}^{m} r_{i}(t) x_{i}\right\|^{2} d t=\sum_{i=1}^{m}\left\|x_{i}\right\|^{2}
$$

valid in suchs spaces) but in general $T_{p}(X)$ may tend to infinity with $\operatorname{dim} X$ if $p>1$. It is known that if an infinite dimensional space $X$ fails to be of type $p$ for any $p>1$ then for
every $\varepsilon>0, X$ has a subspace $Y_{\varepsilon}$ with $d\left(Y_{\varepsilon}, l_{2}\right) \leqslant 1+\varepsilon$. In other words such $X$ contains for every $\varepsilon>0$ unit vectors $u_{\varepsilon}, v_{\varepsilon}$ so that $\left\|u_{\varepsilon}+v_{\varepsilon}\right\| \geqslant 2-\varepsilon$ and $\left\|u_{\varepsilon}-v_{\varepsilon}\right\| \geqslant 2-\varepsilon$ and thus the $\delta$ defined by (1.6) is 0 (cf. [L.Tz.2] or [M.Sche] for further references and other basic facts concerning the notion of type). Hence by a straightforward compactness argument we get that for a Banach space $X$ (finite or infinite dimensional) with a given $\delta>0$ defined by (1.6) there is a $p(\delta)>1$ and a finite $C(\delta)$ so that $T_{p(\delta)}(X) \leqslant C(\delta)$. It is not hard to give a good explicit estimate of $p(\delta)$ and $C(\delta)$ in terms of $\delta$ but this is of no real relevance to the problems we consider in this paper. We shall from now on state our results in terms of $p$ and $T_{p}(X)$ without restating them in terms of $\delta$.

We exploit the uniform-convexity assumption via the next proposition which is a special case of a factorization theorem of Pisier [P.2, Theorem 1.2] (which in turn is based on earlier results of Maurey and Nikishin). In order to state the proposition we have to recall the notion of weak $L_{p}$ spaces. Let $(\Omega, \mu)$ be a probability measure space and let $p \geqslant 1$. We denote by $L_{p, \infty}(\mu)$ the space of real valued measurable functions $f$ on $\Omega$ so that

$$
\begin{equation*}
\|f\|_{p, \infty}=\sup _{>0} t(\mu\{\omega \in \Omega ;|f(\omega)|>t\})^{1 / p}<\infty . \tag{3.2}
\end{equation*}
$$

Proposition 3.1 (Pisier). Let $X$ be a subspace of $L_{1}(\Omega, \mu)$ of type $p$ for some $p>1$. Then there is a non-negative function $F$ in $L_{1}(\Omega, M)$ so that

$$
\|F\|_{L_{1}(\omega)}=1, \quad\{\omega ; F(\omega)=0\} \subset \cap_{x \in X}\{\omega ; x(\omega)=0\}
$$

and so that for every $x \in X$

$$
\begin{equation*}
\|x(\omega) / F(\omega)\|_{L_{p, \alpha}(\Omega, F \mu)} \leqslant e T_{p}(X)\|x\|_{L_{1}(\Omega, \mu)} \tag{3.3}
\end{equation*}
$$

where $e$ is the base of the natural logarithm.

In other words the proposition asserts that after a suitable change of density transformation (passing from $\mu$ to the measure $\nu$ defined by $d \nu=F d \mu$ ) we can assume that the $L_{p, \infty}$ norm on the unit ball of $X$ is bounded by a constant which is essentially $T_{p}(X)$.

It will be very convenient for our purposes to work not with general change of density operations but only with those for which $F(\omega) \geqslant 1 / 2$ for every $\omega \in \Omega$. This can always be achieved at the cost of an increase of the constant by an amount depending on $p$.

Lemma 3.2. Let $X$ be a subspace of $L_{1}(\Omega, \mu)$ of type $p>1$. Then there is a function $G$ in $L_{1}(\Omega, \mu)$ so that

$$
\begin{equation*}
G(\omega) \geqslant 1 / 2 \quad \text { for every } \quad \omega \in \Omega, \quad\|G(\omega)\|_{L_{1}(\mu)}=1 \tag{3.4}
\end{equation*}
$$

and so that for every $x \in X$

$$
\begin{equation*}
\|x(\omega) / G(\omega)\|_{L_{p, \alpha}(\Omega, G \mu)} \leqslant c T_{p}(X)\|x\| / /(p-1) \tag{3.5}
\end{equation*}
$$

where $c$ is an absolute constant.
Proof. Let $F$ be the function given by Proposition 3.1. The function $G=(1+F) / 2$ clearly satisfies (3.4). We shall show that also (3.5) holds. Let $v$ be the measure on $\Omega$ defined by $d v=F d \mu$ let $t>0$ and let $x \in X$ with $\|x\|=1$. Put

$$
\begin{gathered}
A_{0}=\{\omega ; F(\omega) \geqslant 1\}, \quad A_{k}=\left\{\omega ; 2^{-k} \leqslant F(\omega)<2^{-k+1}\right\}, \quad k=1,2, \ldots \\
D=\{\omega ;|x(\omega)| \geqslant t(1+F(\omega)) / 2\}
\end{gathered}
$$

Clearly

$$
D \subset\{w ;|x(\omega)| \geqslant t F(\omega) / 2\} \cap\{\omega ;|x(\omega)| \geqslant t / 2\}
$$

and hence by (3.3)

$$
v(D) \leqslant v\{\omega ;|x(\omega)| \geqslant t F(\omega) / 2\} \leqslant\left(2 e T_{p}(X) / t\right)^{p}
$$

and also

$$
\begin{aligned}
\mu(D) & =\sum_{k=0}^{\infty} \mu\left(D \cap A_{k}\right) \leqslant \sum_{k=0}^{\infty} 2^{k} v\left(D \cap A_{k}\right) \\
& \leqslant v(D)+\sum_{k=1}^{\infty} 2^{k} v\left\{\omega ;|x(\omega)|>t 2^{k-2} F(\omega)\right\} \\
& \leqslant\left(\frac{e T_{p}(X)}{t}\right)^{p}\left(2^{p}+\sum_{k=1}^{\infty} 2^{k} 2^{p(2-k)}\right)
\end{aligned}
$$

Consequently, for some absolute $c_{1}$

$$
\frac{\mu(D)+v(D)}{2} \leqslant\left(\frac{c_{1} T_{p}(X)}{(p-1) t}\right)^{p}
$$

and this is equivalent to (3.5).

The main technical result in this section is the following

Lemma 3.3. Let $X$ be a subspace of $l_{1}^{N}$ of dimension $n$. Let $1<p \leqslant 2$ and $0<\varepsilon<1 / 2$ and assume that $N \geqslant n \varepsilon^{-2}$. Then there is a subspace $Y$ of $\hat{l}_{1}^{\tilde{N}}$ with $d(X, Y) \leqslant 1+\varepsilon$ provided that

$$
\begin{equation*}
\tilde{N} \geqslant c T_{p}(X) N^{1 / p} n^{1-1 / p} \varepsilon^{-2}\left(\log \varepsilon^{-1}\right)^{1-1 / p}\left(\log \frac{N}{n}\right)^{1 / p}(p-1)^{-3} \tag{3.6}
\end{equation*}
$$

where $c$ is an absolute constant.

Proof. By Lemmas 2.3 and 3.2 we may assume that there is a probability measure $\mu$ on $\{1,2, \ldots, N\}$ so that $\mu\{i\} \geqslant 1 / 2 N$ for every $i$ and so that

$$
\begin{equation*}
\|x\|_{L_{\alpha}(\mu)} \leqslant c_{2} n\|x\| ; \quad\|x\|_{L_{p, \infty}(\mu)} \leqslant c_{2} T_{p}(X)\|x\| /(p-1), \quad x \in X . \tag{3.7}
\end{equation*}
$$

By Lemma 2.4 there is an $\varepsilon$-net $\mathscr{F}$ in the unit ball of $X$ so that

$$
\begin{equation*}
\log \stackrel{=}{\mathscr{F}} \leqslant c_{3} n \log \varepsilon^{-1} \tag{3.8}
\end{equation*}
$$

For every $x \in \mathscr{F}$ we put

$$
\begin{equation*}
A_{1}(x)=\{i ;|x(i)| \leqslant 2\}, \quad A_{k}(x)=\left\{i ; 2^{k-1}<|x(i)| \leqslant 2^{k}\right\}, \quad k=2,3, \ldots \tag{3.9}
\end{equation*}
$$

Note that by (3.7) $A_{k}(x)$ is empty for $k>\left[\log c_{2} n\right]$. We also put for $x \in \mathscr{F}, x_{k}=x \cdot \chi_{A_{k}(x)}$, $k=1,2, \ldots$ and note that

$$
\begin{equation*}
x=\sum_{k=1}^{\left[\log c_{2} n\right]} x_{k}, \quad\left|x_{k}\right| \wedge\left|x_{h}\right|=0 \quad \text { if } \quad h \neq k, \quad x \in \mathscr{F} . \tag{3.10}
\end{equation*}
$$

We shall specify below a choice of positive scalars $\varepsilon_{k}, k=1, \ldots$ so that

$$
\begin{equation*}
\sum_{k=1}^{\left[\log c_{2} n\right]} \varepsilon_{k} \leqslant c_{4} \varepsilon /(p-1) \tag{3.11}
\end{equation*}
$$

and prove, applying the empirical distribution method to the families $\left\{x_{k}\right\}_{x \in \mathscr{F}}$, $k=1,2 \ldots$ that if

$$
\begin{equation*}
N_{0} \geqslant c_{5} T_{p}(X) N^{1 / p} n^{1-1 / p} \varepsilon^{-2}\left(\log \varepsilon^{-1}\right)^{1-1 / p}\left(\log \frac{N}{n}\right)^{1 / p}(p-1)^{-1} \tag{3.12}
\end{equation*}
$$

there is a choice of integers $\left\{i_{j}\right\}_{j=1}^{N}$ (not necessarily distinct) in $\{1,2, \ldots, N\}$ so that

$$
\begin{equation*}
\left|\left\|x_{k}\right\|-N_{0}^{-1} \sum_{j=1}^{N_{0}}\right| x_{k}\left(i_{j}\right)| | \leqslant \varepsilon_{k}, \quad x \in \mathscr{F}, \quad k=1,2 \ldots \tag{3.13}
\end{equation*}
$$

Once we verify (3.13) we are done. Indeed by (3.10) and (3.11) it will follow that

$$
\left|\|x\|-N_{0}^{-1} \sum_{j=1}^{N_{0}}\right| x\left(i_{j}\right)\left|\mid \leqslant \sum_{k} \varepsilon_{k} \leqslant c_{4} \varepsilon /(p-1), \quad x \in \mathscr{F}\right.
$$

Consequently, by applying Lemma 2.5 to $T: X \rightarrow l_{1}^{N_{0}}$ defined by $T x=N_{0}^{-1}\left(x\left(i_{1}\right), \ldots, x\left(i_{N_{0}}\right)\right.$, we get that $d(X, T X) \leqslant c_{6} \varepsilon /(p-1)$. We have only now to replace $\varepsilon$ by $c_{6} \varepsilon /(p-1)$; this transforms the $N_{0}$ of (3.12) to the $\tilde{N}$ of (3.6) and will conclude the proof.

In order to verify (3.13) note first that by (3.7) and (3.9) we get for every $x \in \mathscr{F}$ and $k \geqslant 2$ that

$$
\begin{equation*}
\overline{\overline{A_{k}(x)}} \leqslant 2 N \mu\left(A_{k}(x)\right) \leqslant\left[c_{7} N\left(T_{p}(X) /(p-1) 2^{h}\right)^{p}\right]=u_{k} \tag{3.14}
\end{equation*}
$$

Consequently, all the elements in $\left\{\left|x_{k}\right|\right\}_{k \in \mathscr{F}}, k=2, \ldots$ are convex combinations of the family $\mathscr{D}_{k}$ of all the functions on $\{1,2, \ldots, N\}$ which vanish outside a certain subset of cardinality $u_{k}$ and on this set take one of the three values $0,2^{k-1}$ or $2^{k}$. In order to ensure that (3.13) holds for a certain $k$ it suffices that a similar statement holds for all $y \in \mathscr{D}_{k}$. The cardinality of $\mathscr{D}_{k}$ is at most $\binom{N}{u_{k}} 3^{u_{k}}$ and hence

$$
\begin{equation*}
\log \overline{\mathscr{D}}_{k} \leqslant c_{8} u_{k} \log \left(1+\frac{N}{u_{k}}\right) \leqslant c_{9} k N\left(\frac{T_{p}(X)}{p-1}\right)^{p} 2^{-k p} \tag{3.15}
\end{equation*}
$$

Let $k_{0}$ be the smallest integer such that for $k \geqslant k_{0}, \overline{\mathscr{D}}_{k} \leqslant \overline{\mathscr{F}}$. By (3.8) and (3.15) we have

In view of Lemma 2.1 we can ensure that (3.13) holds for all $k$ provided that

$$
\begin{equation*}
\mathscr{F} \sum_{k=1}^{k_{0}} \exp \left(-\varepsilon_{k}^{2} N_{0} / 8 \cdot 2^{k}\right)+\sum_{k=k_{0}+1}^{\left[\log c_{2} n\right]} \overline{\bar{D}}_{k} \exp \left(-\varepsilon_{k}^{2} N_{0} / 8 \cdot 2^{k}\right)<1 / 2 \tag{3.17}
\end{equation*}
$$

We choose first $\varepsilon_{k_{0}}=\varepsilon$ and $\varepsilon_{k_{0}-h}=(2 / 3)^{h / 2} \varepsilon$ for $h=1,2, \ldots, k_{0}-1$. Then $\Sigma_{k=1}^{k_{0}} \varepsilon_{k} \leqslant \varepsilon /(1-\sqrt{2 / 3})$ and so far we fulfil (3.11). If we put $\eta=\exp \left(-\varepsilon^{2} N_{0} / 8 \cdot 2^{k_{0}}\right)$ then in view of (3.8), (3.12) and
(3.16) we get that $\overline{\mathscr{F}} \eta<\frac{1}{8}$ (provided the constant $c_{5}$, which we are still free to choose, is large enough). Hence

$$
\begin{aligned}
\stackrel{\mathscr{F}}{\mathcal{F}_{k}} \sum_{k=1}^{k_{0}} \exp \left(-\varepsilon_{k}^{2} N_{0} / 8 \cdot 2^{k}\right) & \leqslant \overline{\mathscr{F}}\left(\eta+\eta^{4 / 3}+\eta^{(4 / 3)^{2}}+\ldots\right) \\
& \leqslant \overline{\mathscr{F}} \eta\left(1+\eta^{1 / 3}+\eta^{2 / 3}+\ldots\right) \leqslant \overline{\overline{\mathscr{F}}} \eta /\left(1-\eta^{1 / 3}\right) \leqslant 2 \text { 产 } \eta<1 / 4
\end{aligned}
$$

We pass to the second sum in (3.17). We put $k_{1}=\left[\log c_{2} n\right]$. If $k_{1} \leqslant k_{0}$ the sum does not appear and there is no problem. So we assume that $k_{0}<k_{1}$. To ensure that (3.17) holds it suffices to choose $\varepsilon_{k}$ for $k_{0}<k \leqslant k_{1}$ so that

$$
\begin{gather*}
\mathscr{\mathscr { D }}_{k} \leqslant\left(\exp \left(\varepsilon_{k}^{2} N_{0} / 8 \cdot 2^{k}\right)\right)^{1 / 2}, \quad k_{0}<k \leqslant k_{1}  \tag{3.18}\\
\sum_{k=k_{0}+1}^{k_{1}}\left(\exp \left(-\varepsilon_{k}^{2} N_{0} / 8 \cdot 2^{k}\right)\right)^{1 / 2}<\frac{1}{4} \tag{3.19}
\end{gather*}
$$

The condition (3.18) is equivalent to

$$
2 c_{9} k N\left(\frac{T_{p}(X)}{p-1}\right)^{p-1} 2^{-k p} \leqslant \varepsilon_{k}^{2} N_{0} \cdot 2^{-k+3}, \quad k>k_{0}
$$

but since $N_{0}$ was chosen so that this holds for $k=k_{0}$ we get that (3.18) is verified once

$$
\varepsilon_{k_{0}+h} \geqslant \varepsilon\left(\frac{k_{0}+h}{k_{0}} 2^{-h(p-1)}\right)^{1 / 2}, \quad h \geqslant 1 .
$$

As for (3.19) we choose $\varepsilon_{k_{1}}=\varepsilon$. Then

$$
\left(\exp \left(-\varepsilon^{2} N_{0} / 8 \cdot 2^{k_{1}}\right)\right)^{1 / 2}<1 / 8
$$

provided $N_{0}>c_{11} n \varepsilon^{-2}$. This requirement on $N_{0}$ is ensured by (3.12) since we assume from the start that $N \geqslant n \varepsilon^{-2}$. Now the same computation as that done for the first sum in (3.17) shows that if

$$
\begin{equation*}
\varepsilon_{k_{1}-h} \geqslant \varepsilon\left(\frac{2}{3}\right)^{h / 2}, \quad h>0 \tag{3.21}
\end{equation*}
$$

then (3.19) holds. We take now $\varepsilon_{k}$ for $k_{0}<k<k_{1}$ to be the maximum of the two numbers appearing in the right hand sides of (3.20) and (3.21). With this choice of the $\left\{\varepsilon_{k}\right\}_{k=1}^{k_{1}}$ we get that (3.17) holds and that $\Sigma_{k=1}^{k_{1}} \varepsilon_{k}$ is bounded by an absolute constant times $\varepsilon /(p-1)$.

Lemma 3.3 is evidently a useful tool for an iteration procedure. Let $X$ be an $n$-dimensional subspace of $L_{1}(0,1)$ and let $\delta_{1}>0$. We find a subspace $Y_{1}$ in $l_{1}^{N_{1}}$ with $d\left(X, Y_{1}\right) \leqslant 1+\delta_{1}$ so that $N_{1}$ is not too large (say by using Theorem 2.6). We then choose another number $\delta_{2}$ and by applying Lemma 3.3 find a subspace $Y_{2}$ of $l_{1}^{N_{2}}$ (with $N_{2}$ the $\tilde{N}$ of Lemma 3.3) with $d\left(Y_{1}, Y_{2}\right) \leqslant 1+\delta_{2}$. We continue to apply Lemma 3.3 as long as we can and as long as the numbers $N_{k}$ decrease. If we stop after $s$ stages we get a space $Y_{s}$ for which $d\left(X, Y_{s}\right) \leqslant \prod_{k=1}^{s}\left(1+\delta_{k}\right)$.

To put this procedure into a precise framework we recall the quantity $N(X, \varepsilon)$ defined by the formulation of problem $\left({ }^{*}\right)$ in the introduction. Since

$$
(1+\varepsilon)\left(1+\frac{\delta}{2}\right) \leqslant 1+\varepsilon+\delta \quad \text { if } \quad 0<\varepsilon, \delta \leqslant 1 / 2
$$

we can restate Lemma 3.3 as follows.

Lemma 3.3'. There is an absolute constant $c$ so that for every choice of $0<\varepsilon$, $\delta<1 / 2$, every $1<p \leqslant 2$ and every $n$-dimensional subspace $X$ of $L_{1}(0,1)$

$$
\begin{equation*}
N(X, \varepsilon+\delta) \leqslant c T_{p}(X) N(X, \varepsilon)^{1 / p} n^{1-1 / p} \delta^{-2}\left(\log \delta^{-1}\right)^{1-1 / p^{\prime}}\left(\log \frac{N(X, \varepsilon)}{n}\right)^{1 / p}(p-1)^{-3} \tag{3.22}
\end{equation*}
$$

provided that $N(X, \varepsilon) \geqslant 4 n \delta^{-2}$.

Corollary 3.4. Let $X$ be an $n$-dimensional subspace of $L_{1}(0,1)$, let $1<p \leqslant 2$, $0<\varepsilon<1$ and let $\varrho>2$. Then

$$
\begin{equation*}
N(X, \varepsilon) \leqslant C\left(p, T_{p}(X), \varrho\right) \varepsilon^{-\varrho p /(p-1)} n \tag{3.23}
\end{equation*}
$$

where, as the notation indicates, $C\left(p, T_{p}(X), \varrho\right)$ depends just on $p, \varrho$ and the type $p$ constant of $X$.

Proof. Let $\tau>1$ be such that $\varrho p /(p-1)>2 \tau p /(p-\tau)$ and put $\varphi(\varepsilon)=N(X, \varepsilon) / n$. By (3.22) there is a $\gamma=\gamma\left(p, T_{p}(X), \tau\right)$ so that for every $0<\varepsilon<1$

$$
\begin{equation*}
\varphi(\varepsilon) \leqslant \gamma \cdot \varepsilon^{-2 \tau} \varphi\left(\frac{\varepsilon}{2}\right)^{\tau / p} \tag{3.24}
\end{equation*}
$$

provided $\varphi(\varepsilon / 2) \geqslant(4 / \varepsilon)^{2}$. By iterating (3.24) $k$ times we get

$$
\begin{equation*}
\varphi(\varepsilon) \leqslant \gamma^{1+\tau / p+\ldots+(\tau / p)^{k}} \varepsilon^{-2 \tau\left(1+\tau / p+\ldots+(\tau / p)^{k}\right)} 2^{\tau / p+2 \tau^{2} / p^{2}+\ldots+k(\tau / p)^{k}} \varphi\left(\frac{\varepsilon}{2^{k}}\right)^{\tau^{k} / p^{k}} \tag{3.25}
\end{equation*}
$$

By Theorem 2.6, $\varphi\left(\varepsilon / 2^{k}\right) \leqslant n \varepsilon^{-2 \tau} 2^{2 k \tau}$. Substituting this inequality in (3.25) with $k \sim \log n$ we get the desired result. (If the iteration stops earlier, i.e., $\varphi\left(\varepsilon / 2^{i}\right) \leqslant\left(2^{i+1} / \varepsilon\right)^{2}$ for some $i<k$ then we have an even better inequality.)

Corollary 3.4 is a weak version of Theorem 1.1 and as mentioned in the introduction it already generalizes the main result of [J.Sche.1]. In order to obtain Theorem 1.1 as stated we have to add to the method of this section an argument involving certain entropy numbers, which will be evaluated in the next section.

## 4. Entropy estimates

In this section we obtain an estimate on the number of balls of radius $t$ in $\left\|\|_{\infty}\right.$ norm needed to cover the unit ball of a subspace $X$ of $L_{1}$ which is in a "good position". The estimate is based on several known results. We devote a large part of this section to the statement of those known (or essentially known) results and the explanation of the notions from Banach space theory which enter in their formulations.

We start with the notion of entropy. Let $D$ and $B$ be subsets of a linear space and let $t>0$. We put

$$
E(D, B, t)=\min \left\{k ; \exists\left\{x_{i}\right\}_{i=1}^{k}, D \subset \bigcup_{i=1}^{k}\left(x_{i}+t B\right)\right\}
$$

In the cases we consider here this number will always be finite. Sometimes it is convenient to use instead of $E(D, B, t)$ the number

$$
\tilde{E}(D, B, t)=\min \left\{k ; \exists\left\{x_{i}\right\}_{i=1}^{k}, x_{i} \in D, D \subset \bigcup_{i=1}^{k}\left(x_{i}+t B\right)\right\}
$$

It follows directly from the definition that always (if $B$ is convex and symmetric)

$$
\begin{equation*}
\tilde{E}(D, B, 2 t) \leqslant E(D, B, t) \leqslant \tilde{E}(D, B, t) \tag{4.1}
\end{equation*}
$$

An important quantity which enters into entropy computations for subsets of $\mathbf{R}^{n}$ is the following: Let $\left|\left|\left|\left|\left|\mid\right.\right.\right.\right.\right.$ denote the Euclidean norm in $\mathbf{R}^{n}$ (with $B^{n}$ its unit ball and $S^{n-1}$ its boundary) and let $\left\|\|\right.$ be another norm on $\mathbf{R}^{n}$. We denote $X=\left(\mathbf{R}^{n},\| \|\right)$ and $X^{*}=\left(\mathbf{R}^{n},\| \|_{*}\right)$ where the duality is with respect to the inner product defined by ||| |||. The average of $\|x\|$ on $S^{n-1}$ is denoted by $M_{X}$, i.e.,

$$
\begin{equation*}
M_{X}=\int_{S^{n-1}}\|x\| d \sigma(x) \tag{4.2}
\end{equation*}
$$

where $\sigma$ is the normalized rotation invariant measure on $S^{n-1}$. By the homogeneity of $\|\cdot\|$ and by using polar coordinates we get easily other convenient formulas for $M_{X}$ for example

$$
\begin{equation*}
M_{X}=\alpha_{n}(2 \pi)^{-n / 2} \int_{\mathbf{R}^{n}}\|x\| \exp \left(-\|x\|^{2} / 2\right) d x, \quad \alpha_{n} \sim n^{-1 / 2} \tag{4.3}
\end{equation*}
$$

A probabilistic way to rewrite (4.3) is to consider $n$ independent and normalized Gaussian variables $\left\{g_{i}(\omega)\right\}_{i=1}^{n}$ on some probability measure space $(\Omega, \mu)$. Then

$$
\begin{equation*}
M_{X}=\alpha_{n} \int_{\Omega}\left\|\sum_{i=1}^{n} g_{i}(\omega) e_{i}\right\| d \mu(\omega) \tag{4.4}
\end{equation*}
$$

where $\left\{e_{i}\right\}_{i=1}^{n}$ is an orthonormal basis in $\mathbf{R}^{n}$.
The following result is well known in probability theory.

Proposition 4.1 (Sudakov [Su]). Let $X=\left(\mathbf{R}^{n}, \|| |\right)$ and let ||| ||| be the Euclidean norm on $\mathbf{R}^{n}$. Then

$$
\begin{equation*}
\log E\left(B_{X}, B^{n}, t\right) \leqslant C n\left(\frac{M_{X^{*}}}{t}\right)^{2} \tag{4.5}
\end{equation*}
$$

where $B_{X}\left(\right.$ resp. $\left.B^{n}\right)$ denote the unit balls of $|||\mid($ resp. $| || || |)$ and $C$ is a universal constant.
A dual version of this proposition was recently obtained by A. Pajor and N. Tomczak-Jaegermann [Pa.T-J].

Proposition 4.2. With the same notation as in Proposition 4.1 we have

$$
\begin{equation*}
\log E\left(B^{n}, B_{X}, t\right) \leqslant c n\left(\frac{M_{X}}{t}\right)^{2} \tag{4.6}
\end{equation*}
$$

We shall present here a proof of (4.6) which is different from that indicated in [Pa.T-J]. This proof was kindly communicated to us by A. Pajor and is due to him and M. Talagrand.

Proof. We let $\mu$ be the probability measure on $\mathbf{R}^{n}$ defined by

$$
d \mu=(2 \pi)^{-n / 2} \exp \left(-\| \| x\| \|^{2} / 2\right) d x
$$

By (4.3)

$$
\begin{equation*}
\mu\left\{x ;\|x\| \leqslant 2 M_{X} \alpha_{n}^{-1}\right\} \geqslant 1 / 2 \tag{4.7}
\end{equation*}
$$

Let $\left\{x_{i}\right\}_{i=1}^{N}$ be a maximal subset in $B^{n}$ relative to the requirement that $\left\|x_{i}-x_{j}\right\| \geqslant t$ for $i \neq j$. The sets $\left\{x_{i}+\frac{1}{2} t B_{X}\right\}_{i=1}^{N}$ have mutually disjoint interiors. Hence,

$$
\begin{equation*}
1 \geqslant \sum_{i=1}^{N} \mu\left\{y_{i}+2 M_{X} \alpha_{n}^{-1} B_{X}\right\} \quad \text { where } \quad y_{i}=4 M_{X}\left(t \alpha_{n}\right)^{-1} x_{i} \tag{4.8}
\end{equation*}
$$

Fix $1 \leqslant i \leqslant N$. By (4.7), the convexity of the function $e^{-u}$ and the symmetry of $B_{X}$ with respect to the origin, we get

$$
\begin{aligned}
\mu\left\{y_{i}+2 M_{X} a_{n}^{-1} B_{X}\right\} & =(2 \pi)^{-n / 2} \int_{2 M_{X} a_{n}^{-1} B_{X}} \exp \left(-\| \| x-y_{i}\| \|^{2} / 2\right) d x \\
& \geqslant(2 \pi)^{-n / 2} \int_{2 M_{X} a_{n}^{-1} B_{X}} \exp \left(-\left(\left|\left\|x-y_{i}\right\|^{2}+\left\|\mid x+y_{i}\right\| \|^{2}\right) / 4\right) d x\right. \\
& =(2 \pi)^{-n / 2} \int_{2 M_{X} a_{n}^{-2} B_{X}} \exp \left(-\left(\||x|\|^{2}+\left\|\mid y_{i}\right\| \|^{2}\right) / 2\right) d x \\
& \geqslant \frac{1}{2} \exp \left(-\|\mid\| y_{i}\| \|^{2} / 2\right) \geqslant \frac{1}{2} \exp \left(-\left(4 M_{X} / \alpha_{n} t\right)^{2}\right)
\end{aligned}
$$

Hence by (4.8), $N \leqslant 2 \exp \left(4 M_{X} / \alpha_{n} t\right)^{2}$ and since $\alpha_{n} \sim n^{-1 / 2}$ we deduce (4.6).

We recall now the notion of the $K$-convexity constant $K(X)$ of the Banach Space $X$. Let $Q$ be the natural projection from $L_{2}([0,1], X)$ onto the span of $\left\{x_{i} r_{i}(t)\right\}_{i=1}^{\infty}$ where $x_{i} \in X$ and $r_{i}(t)$ are the Radamacher functions i.e.,

$$
\begin{equation*}
Q f(s)=\sum_{i=1}^{\infty}\left(\int_{0}^{1} f(t) r_{i}(t) d t\right) r_{i}(s), \quad f \in L_{2}([0,1], X) \tag{4.9}
\end{equation*}
$$

Then

$$
\begin{equation*}
K(X)=\|Q\| \quad \text { (with the usual operator norm) } \tag{4.10}
\end{equation*}
$$

It is evident that $\|Q\|=1$ if $X$ is a Hilbert space and thus in particular $K(X)<\infty$ whenever $\operatorname{dim} X<\infty$. The $K$ convexity constant will enter our computations in view of the
following inequality which is an immediate consequence of the definition. For every choice of $\left\{x_{i}^{*}\right\}_{i=1}^{m}$ in $X^{*}$ we have

$$
\begin{align*}
& \left(\int_{0}^{1} \| \sum_{i=1}^{m} r_{i}(t) x_{i}^{*}| |_{X^{*}}^{2} d t\right)^{1 / 2}=\sup \left\{\int_{0}^{1}\left\langle\sum_{i=1}^{m} r_{i}(t) x_{i}^{*}, f(t)\right\rangle d t,\|f\|_{L_{2}(X)} \leqslant 1\right\} \\
& \quad \leqslant K(X) \sup \left\{\int_{0}^{1}\left\langle\sum_{i=1}^{m} r_{i}(t) x_{i}^{*}, \sum_{i=1}^{m} r_{i}(t) x_{i}\right\rangle d t,\left.\left\|\sum_{i=1}^{m} r_{i}(t) x_{i}\right\|\right|_{L_{2}(X)} \leqslant 1\right\}  \tag{4.11}\\
& \quad=K(X) \sup \left\{\sum_{i=1}^{m} x_{i}^{*}\left(x_{i}\right) ; \| \sum_{i=1}^{m} r_{i}(t) x_{i}\right\} \left\lvert\,\left\{\begin{array}{|l} 
\\
L_{2}(x)
\end{array} \leqslant\right.\right.
\end{align*}
$$

As for estimates on $K(X)$; the following are due to Pisier [P.3] (see also [M.Sche], Chapter 14).

$$
\begin{gather*}
K(X) \leqslant C \log d\left(X, l_{2}^{n}\right) \leqslant C \log n, \quad n=\operatorname{dim} X  \tag{4.12}\\
K(X) \leqslant C(\log n)^{1 / 2}, \quad n=\operatorname{dim} X, \quad X \subset L_{1}(0,1)  \tag{4.13}\\
K\left(L_{p}(\Omega, \mu)\right) \sim(p-1)^{1 / 2}, \quad 1<p \leqslant 2 \tag{4.14}
\end{gather*}
$$

An estimate which seems to be new is given in

Lemma 4.3. Let $X$ be an $n$-dimensional subspace of $L_{1}(0,1)$ and let $1<p \leqslant 2$. Then

$$
\begin{equation*}
K(X) \leqslant C\left(\log \left(1+T_{p}(X)\right) /(p-1)\right)^{1 / 2} \tag{4.15}
\end{equation*}
$$

where $C$ is an absolute constant and $T_{p}(X)$ the type $p$ constant of $X$.

Proof. By Proposition 3.1 there is no loss of generality to assume that $X \subset L_{1}(\Omega, \mu)$ for some probability measure space $(\Omega, \mu)$ and that we have for every $x \in X$

$$
\begin{equation*}
\|x\|_{L_{p, \infty}} \leqslant e T_{p}(X)\|x\|_{L_{1}} \tag{4.16}
\end{equation*}
$$

Let $0<\theta<1$ and $1 / s=(1-\theta) / 1+\theta / p$. A standard computation (involving the sets $\left\{\omega, 2^{k} \leqslant|f(\omega)| \leqslant 2^{k+1}\right\}$ for every $k$ ) shows that for $f \in L_{p, \infty}(\Omega, \mu)$ we have for some absolute $c_{1}$

$$
\|f\|_{L_{s}} \leqslant c_{1}\left(\frac{p-1}{p-s}\right)^{1 / s}\|f\|_{L_{1}}^{1-\theta}\|f\|_{L_{p, x}}^{\theta}
$$

and in particular if we restrict $\theta$ to be in $\left(0, \frac{1}{2}\right]$

$$
\begin{equation*}
\|f\|_{L_{s}} \leqslant \varepsilon_{2}\|f\|_{L_{1}}^{1-\theta}\|f\|_{L_{p, \infty}}^{\theta} \tag{4.17}
\end{equation*}
$$

for some absolute $c_{2}$. Hence by (4.16) and (4.17) we get that if $\theta \leqslant \frac{1}{2}$ then the Banach Mazur distance of $X$ from a subspace, of $L_{s}(\Omega, \mu)$ is $\leqslant c_{3} T_{p}(X)^{\theta}$. Consequently by (4.14) $K(X) \leqslant c_{4}(s-1)^{-1 / 2} T_{p}(X)^{\theta}$. By choosing $\theta$ so that $2(s-1) \log \left(1+T_{p}(X)\right)=p-1$ we get (4.15).

Two other notions from Banach space theory will be useful for us: 2-absolutely summing norms and cotype $q$. The 2 -absolutely summing norm $\pi_{2}(T)$ of a linear operator $T: X \rightarrow Y$ is defined to be the smallest constant $\gamma$ so that

$$
\begin{equation*}
\sum_{i=1}^{m}\left\|T x_{i}\right\|^{2} \leqslant \gamma^{2} \sup \left\{\sum_{i=1}^{m}\left|x^{*}\left(x_{i}\right)\right|^{2},\left\|x^{*}\right\|_{X^{*}} \leqslant 1\right\} \tag{4.18}
\end{equation*}
$$

holds for every choice of $\left\{x_{i}\right\}_{i=1}^{m}, m=1,2, \ldots$, in $X$ (in our case $\operatorname{dim} X$ will be finite and this ensures in particular that $\pi_{2}(T)<\infty$ for every operator $T$ ). The notion of cotype is dual to that of type. For $2 \leqslant q<\infty$ the cotype $q$ constant of the Banach space $Y$, denoted by $C_{q}(Y)$ is the smallest constant $\gamma$ for which

$$
\begin{equation*}
\left(\sum_{i=1}^{m}\left\|y_{i}\right\|^{q}\right)^{1 / q} \leqslant \gamma \int_{0}^{1}\left\|\sum_{i=1}^{m} y_{i} r_{i}(t)\right\| d t \tag{4.19}
\end{equation*}
$$

holds for every choice of $\left\{y_{i}\right\}_{i=1}^{m}, m=1,2, \ldots$, in $Y$. A simple consequence of the classical Khintchine inequality is that $C_{2}\left(L_{1}(0,1)\right)<\infty$, i.e., $L_{1}(0,1)$ is a space of cotype 2. With these notions we can state the next lemma which is just a reformulation of Lemma 1 in [D.M.T-J].

Lemma 4.4. With $X$ as in Proposition 4.1. we have

$$
\begin{equation*}
M_{X^{*}} \leqslant c n^{-1 / 2} K(X) C_{2}(X) \pi_{2}(T) \tag{4.20}
\end{equation*}
$$

where $T$ is the formal identity map from $X$ into $l_{2}^{n}=\left(\mathbf{R}^{n},|||| |)\right.$. In particular, if $X$ is a subspace of $L_{1}(0,1)$ then

$$
\begin{equation*}
M_{X^{*}} \leqslant \tilde{c} n^{-1 / 2} K(X) \pi_{2}(T) \tag{4.21}
\end{equation*}
$$

Proof. Evidently (4.21) follows from (4.20) with $\tilde{c}=c C_{2}\left(L_{1}(0,1)\right)$. We prove (4.20). By the Pietch factorization theorem (cf. e.g., [L.Tz. 1 p. 64]) there exists an $m \geqslant n$ so that $T$ can be factored as $T=V \Delta U$

$$
X \xrightarrow{U} l_{\infty}^{m} \xrightarrow{\Delta} l_{2}^{m} \xrightarrow{V} l_{2}^{n}
$$

where $\|U\| \leqslant 1,\|V\| \leqslant 1, \Delta f_{j}=\lambda_{j} f_{j}$, the $\left\{f_{j}\right\}_{j=1}^{m}$ being the unit vectors in $l_{\infty}^{m}$ and $l_{2}^{m}$ and

$$
\begin{equation*}
\left(\sum_{j=1}^{m} \lambda_{j}^{2}\right)^{1 / 2}=\|\Delta\| \leqslant 2 \pi_{2}(T) \tag{4.22}
\end{equation*}
$$

Since an orthogonal $m \times m$ matrix takes $m$ independent normalized Gaussian variables $\left\{g_{i}\right\}_{i=1}^{m}$ (on some probability space $(\Omega, \mu)$ ) into a set of $m$ variables having the same joint distribution and since every operator of norm $\leqslant 1$ on $l_{2}^{n}$ is a convex combination of orthogonal transformations we deduce (using (4.4) and the fact that $\left\|V^{*}\right\| \leqslant 1$ ) that

$$
M_{X^{*}}=\alpha_{n} \int_{\Omega}\left\|\sum_{i=1}^{n} g_{i}(\omega) e_{i}\right\|_{X^{*}} d \mu(\omega) \leqslant \alpha_{n} \int_{\Omega}\left\|\sum_{j=1}^{m} g_{j}(\omega) U^{*} \Delta^{*} f_{j}\right\|_{X^{*}} d \mu(\omega)
$$

Again, since for every $t$ the functions $\left\{r_{j}(t) g_{j}(\omega)\right\}_{j=1}^{m}$ have the same joint distribution as $\left\{g_{j}(\omega)\right\}_{j=1}^{m}$ it follows that

$$
\begin{align*}
M_{X^{*}} & \leqslant \alpha_{n} \int_{\Omega} \int_{0}^{1}\left\|\sum_{j=1}^{m} r_{j}(t) g_{j}(\omega) U^{*} \Delta^{*} f_{j}\right\| d t d \mu(\omega)  \tag{4.23}\\
& \leqslant \alpha_{n} \int_{\Omega}\left(\int_{0}^{1}\left\|\sum_{j=1}^{m} g_{j}(\omega) r_{j}(t) U^{*} \Delta^{*} f_{j}\right\|^{2} d t\right)^{1 / 2} d \mu(\omega) .
\end{align*}
$$

For every $\omega \in \Omega$ we get by (4.11)

$$
\begin{gather*}
\left(\int_{0}^{1}\left\|\sum_{j=1}^{m} g_{j}(\omega) r_{j}(t) U^{*} \Delta^{*} f_{j}\right\|^{2}\right)^{1 / 2}  \tag{4.24}\\
\leqslant K(X) \sup \left\{\sum_{j=1}^{m} g_{j}(\omega) U^{*} \Delta^{*} f_{j}\left(x_{j}\right) ;\left\|\sum_{j=1}^{m} r_{j}(t) x_{j}\right\|_{L_{2}(X)} \leqslant 1\right\} .
\end{gather*}
$$

By (4.19) we have for all $\left\{x_{j}\right\}_{j=1}^{m}$ in $X$

$$
\left(\sum_{j=1}^{m}\left\|x_{j}\right\|^{2}\right)^{1 / 2} \leqslant C_{2}(X)\left\|\sum_{j=1}^{m} r_{j}(t) x_{j}\right\|_{L_{1}(X)} \leqslant C_{2}(X)\left\|\sum_{j=1}^{m} r_{j}(t) x_{j}\right\|_{L_{2}(X)}
$$

Hence $\left\|\sum_{j=1}^{m} r_{j}(t) x_{j}\right\|_{L_{2}(X)} \leqslant 1$ implies that

$$
\sum_{j=1}^{m} g_{j}(\omega) U^{*} \Delta^{*} f_{j}\left(x_{j}\right) \leqslant\left(\sum_{j=1}^{m}\left|g_{j}(\omega) \lambda_{j}\right|^{2}\right)^{1 / 2}\left(\sum_{j=1}^{m} \|\left. U x_{j}\right|^{2}\right)^{1 / 2} \leqslant C_{2}(X)\left(\sum_{j=1}^{m}\left|g_{j}(\omega) \lambda_{j}\right|^{2}\right)^{1 / 2}
$$

Combining this inequality with (4.22), (4.23) and (4.24) we get

$$
\begin{aligned}
M_{X^{*}} & \leqslant \alpha_{n} K(X) C_{2}(X) \int_{\Omega}\left(\sum_{j=1}^{m}\left|g_{j}(\omega) \lambda_{j}\right|^{2}\right)^{1 / 2} d \mu(\omega) \\
& \leqslant c_{1} \alpha_{n} K(X) C_{2}(X)\left(\sum_{j=1}^{m} \lambda_{j}^{2}\right)^{1 / 2} \\
& \leqslant c n^{-1 / 2} K(X) C_{2}(X) \pi_{2}(T)
\end{aligned}
$$

Our next lemma is another change of density result.

Lemma 4.5. Let $(\Omega, \mu)$ be a probability space and let $X$ be an $n$ dimensional subspace of $L_{1}(\Omega, \mu)$. Then there is another probability measure $v$ on $\Omega$ and a subspace $\tilde{X}$ of $L_{1}(\Omega, v)$ which is isometric to $X$ so that

$$
\begin{equation*}
\pi_{2}(T) \leqslant c n^{1 / 2} \tag{4.25}
\end{equation*}
$$

where $T$ is the identity map from $\tilde{X}$ onto $(\tilde{X},|||| |)$ where $||| || |$ is the norm induced on $\tilde{X}$ from $L_{2}(\Omega, v)$ and $c$, as usual, denotes an absolute constant.

Proof. As shown by Lewis [Le] there is a basis $\left\{\varphi_{i}\right\}_{i=1}^{n}$ of $X$ so that for all scalars $\left\{\lambda_{j}\right\}_{j=1}^{n}$

$$
\sum_{j=1}^{n} \lambda_{j}^{2} / n=\int_{\Omega}\left(\left|\sum_{j=1}^{n} \lambda_{j} \varphi_{j}\right|^{2} / \Phi\right) d \mu
$$

where

$$
\Phi=\left(\sum_{j=1}^{n} \varphi_{j}^{2}\right)^{1 / 2}, \quad\|\Phi\|_{L_{1}(\mu)}=1
$$

Define now a measure $v$ on $\Omega$ by $d \nu=\Phi d \mu$ and let $f \rightarrow \Phi^{-1} f$ map $X$ onto a subspace $\tilde{X}$ of $L_{1}(\Omega, \nu)$. Clearly this map is an isometry and $\tilde{X}$ has a basis $\psi_{j}=\varphi_{j} / \Phi, j=1,2, \ldots, n$ satisfying

$$
\begin{equation*}
\sum_{j=1}^{n} \psi_{j}^{2} \equiv 1, \quad \int_{\Omega} \psi_{j}^{2} d \nu=n^{-1}, \quad \int_{\Omega} \psi_{k} \psi_{j} d \nu=0 \quad \text { if } \quad k \neq j \tag{4.26}
\end{equation*}
$$

Consider now the functions $h_{t}$ defined on $\Omega, 0 \leqslant t \leqslant 1$ by $\left(i^{2}=-1\right)$

$$
h_{t}(\omega)=\operatorname{Imag} \prod_{j=1}^{n}\left(1+i r_{j}(t) \psi_{j}(\omega)\right)=\sum_{j=1}^{n} r_{j}(t) \psi_{j}(\omega)+\text { higher order Walsh function in } t
$$

By (4.26)

$$
\begin{equation*}
\left\|h_{t}\right\|_{L_{\infty}(v)} \leqslant \sup _{\omega} \prod_{j=1}^{n}\left(1+\psi_{j}^{2}(\omega)\right)^{1 / 2} \leqslant \sqrt{e} \tag{4.27}
\end{equation*}
$$

and for every $x=\sum_{j=1}^{n} \lambda_{j} \psi_{j}$ in $\tilde{X}$

$$
\int_{\Omega} x(\omega) h_{t}(\omega) d v=\sum_{j=1}^{n} \lambda_{j} r_{j}(t) / n+\text { higher order Walsh functions }
$$

and thus

$$
\begin{equation*}
\int_{0}^{1}\left(\int_{\Omega} x(\omega) h_{t}(\omega) d v\right)^{2} d t \geqslant \sum_{j=1}^{n} \lambda_{j}^{2} / n^{2}=\| \| x \|_{L_{2}(v)}^{2} / n . \tag{4.28}
\end{equation*}
$$

By (4.27) and (4.28) we get for any choice of $\left\{x_{k}\right\}_{k=1}^{m}$ in $X$

$$
\begin{aligned}
\sup \left\{\sum_{k=1}^{m}\left|x^{*}\left(x_{k}\right)\right|^{2} ;\left\|x^{*}\right\| \leqslant 1\right\} & \geqslant e^{-1} \sup _{0 \leqslant t \leqslant 1} \sum_{k=1}^{m}\left(\int_{\Omega} h_{t}(\omega) x_{k}(\omega) d v\right)^{2} \\
& \geqslant e^{-1} \sum_{k=1}^{m} \int_{0}^{1}\left(\int_{\Omega} h_{t}(\omega) x_{k}(\omega) d v\right)^{2} d t \geqslant \sum_{k=1}^{m}\left\|x_{k}\right\| \|^{2} / n e
\end{aligned}
$$

and this is equivalent to (4.25).
Remark. Lemma 4.5 is already the third change of density result mentioned in this paper. By forming averages of the densities obtained in the various results and allowing a possible change in the absolute constants involved, we can ensure that we work with densities which have all the desired properties at once. In all the situations appearing here it is easy to check that the averaging procedure is permissible (the least obvious case involves Proposition 3.1 where the checking is made in Lemma 3.2).

We combine now most of the preceding results to obtain the following entropy estimate.

Proposition 4.6. Let $X$ be an $n$-dimensional subspace of $L_{1}(\Omega, \mu)$. Then there is $a$ probability measure $v$ on $\Omega$ given by $d v=F d \mu$ so that for the space $\tilde{X}=F^{-1} X$ in $L_{1}(\Omega, v)$ we have for every $t>0$ and every $q \geqslant 1$

$$
\begin{equation*}
\log E\left\{B_{\bar{X}}, B_{q}, t\right\} \leqslant c K(X) n q^{1 / 2} / t \tag{4.29}
\end{equation*}
$$

where $c$ is an absolute constant and $B_{q}$ is the unit ball of $\tilde{X}_{q}=\tilde{X}$ endowed with the norm induced by $L_{q}(\Omega, v)$.

Proof. Let $v$ be as in the proof of Lemma 4.5. From the definition of entropy it follows immediately that

$$
\begin{equation*}
E\left\{B_{\tilde{X}}, B_{q}, t\right\} \leqslant E\left\{B_{\tilde{X}}, B_{2}, \lambda\right\} E\left\{B_{2}, B_{q}, t / \lambda\right\} \tag{4.30}
\end{equation*}
$$

for every $0<\lambda<t$. We take as ||| ||| the norm in $\tilde{X}_{2}$ i.e., that induced by $L_{2}(\Omega, v)$. By Propositions 4.1, 4.2 and (4.30) we deduce that

$$
\begin{equation*}
\log E\left\{B_{\bar{X}}, B_{q}, t\right\} \leqslant c_{1} n\left(\frac{M_{X^{*}}^{2}}{\lambda^{2}}+\frac{\lambda^{2}}{t^{2}} M_{X_{q}}^{2}\right) . \tag{4.31}
\end{equation*}
$$

By Lemmas 4.4 and 4.5, $M_{X^{*}} \leqslant c_{2} K(X)$. Let $\left\{g_{j}\left(\omega^{\prime}\right)\right\}_{j=1}^{n}$ be normalized independent Gaussians on a probability space ( $\Omega^{\prime}, \mu^{\prime}$ ). We take as unit vectors in $\tilde{X}_{2}$ the functions $\left\{\sqrt{n} \psi_{j}\right\}_{j=1}^{n}$ with the $\psi_{j}$ as in (4.26). By (4.4)

$$
M_{\tilde{X}_{q}}=\alpha_{n} n^{1 / 2} \int_{\Omega}\left\|\sum_{j=1}^{n} g_{j}\left(\omega^{\prime}\right) \psi_{j}\right\|_{L_{q}\left(\Omega^{\prime}, \nu\right)} d \mu\left(\omega^{\prime}\right) .
$$

It follows from the Maurey-Khintchine inequality (cf. e.g., [L.Tz. 2 p. 49]) (4.26) and the fact that $\alpha_{n} \sim n^{-1 / 2}$ that

$$
M_{\tilde{X}_{q}} \leqslant c_{3} q^{1 / 2}\left\|\left(\sum_{j=1}^{n} \psi_{j}^{2}\right)^{1 / 2}\right\|_{L_{q}(\Omega, v)}=c_{3} q^{1 / 2} .
$$

Substituting in (4.31) we get that

$$
\log E\left\{B_{\tilde{X}}, B_{q}, t\right\} \leqslant c_{4} n\left(\frac{K(X)^{2}}{\lambda^{2}}+\frac{q \lambda^{2}}{t^{2}}\right) .
$$

By taking minimum over $\lambda$, (4.29) follows.

Remark. As in the remark above we may replace $\nu$ by $(\nu+\mu) / 2$. Then we get a measure on $\Omega$ which is dominated from below by $\mu / 2$ and for which (4.29) still holds (with $c$ replaced by $2 c$ ).

Corollary 4.7. Let $X$ be an n-dimensional subspace of $l_{1}^{N}$. Then there exists a probability measure $v$ on $\{1, \ldots, N\}$ and a subspace $\tilde{X}$ of $L_{1}(N, v)$ isometric to $X$ so that

$$
\begin{equation*}
\log E\left\{B_{\tilde{X}}, B_{\infty}, t\right\} \leqslant c(\log n \log N)^{1 / 2} n / t \tag{4.32}
\end{equation*}
$$

where $B_{\infty}$ is the unit ball of $l_{\infty}^{N}=(N, v)$.

Proof. Let $\mu$ be the usual probability measure on $\{1,2, \ldots, N\}$, i.e., $\mu\{i\}=N^{-1}$ for every $i$. We apply Proposition 4.6 to find $\bar{X}$ and $\nu$. By the remark after Proposition 4.6 we may assume that $v\{i\}>(2 N)^{-1}$ for every $i$. For every $f \in L_{q}(N, v)$

$$
\max _{i}\left|f_{i}\right| /(2 N)^{1 / q} \leqslant\|f\|=\left(\sum_{i=1}^{N}\left|f_{i}\right|^{q} v\{i\}\right)^{1 / q} \leqslant \max _{i}\left|f_{i}\right|
$$

and hence $d\left(L_{q}(N, v), L_{\infty}(N, v)\right) \leqslant(2 N)^{1 / q}$. We take now $q=\log N$ in (4.29), use the estimate (4.13) for $K(X)$ and (4.32) follows.

Corollary 4.8. Let $X$ be an $n$-dimensional subspace of $L_{1}(0,1)$. Then for every $\varepsilon>0$ there is a probability space $(\Omega, v)$ and a subspace $\bar{X}$ of $L_{1}(\Omega, v)$ with $d(X, \bar{X}) \leqslant 1+\varepsilon$ so that

$$
\begin{equation*}
\log E\left\{B_{\hat{X}}, B_{L_{\infty}(\Omega, v}, t\right\} \leqslant c\left(\log n\left(\log n \varepsilon^{-1}\right)\right)^{1 / 2} n / t \tag{4.33}
\end{equation*}
$$

Proof. Combine (4.32) with Theorem 2.6.
In the two corollaries above we may use the estimate (4.15) for $K(X)$ instead of (4.13). Then for example we get, instead of (4.33), the following estimate with $1<p \leqslant 2$,

$$
\begin{equation*}
\log E\left\{B_{\tilde{X}}, B_{L_{\infty}(\Omega, v)}, t\right\} \leqslant c\left[\log \left(T_{p}(X)+1\right) \log \left(n \varepsilon^{-1}\right)(p-1)^{-1}\right]^{1 / 2} n / t \tag{4.34}
\end{equation*}
$$

In section 9 we shall prove a slightly stronger version of (4.33), which however will not be of use for estimating $N(X, \varepsilon)$.

## 5. Proof of the main theorems

Our purpose now is to combine the tools developed in the last three sections and to derive from them Theorems 1.1 and 1.2. We show first how the entropy numbers enter into the analysis of the empirical method.

Proposition 5.1. Let $X$ be an n-dimensional subspace of $L_{1}(\Omega, \mu)$ for some probability measure $\mu$. Let $\|\cdot\|_{\infty}$ denote the norm in $L_{\infty}(\Omega, \mu)$ and put

$$
B_{X}=\{x ; x \in X,\|x\| \leqslant 1\}, \quad B_{\infty}=\left\{x ; x \in X,\|x\|_{\infty} \leqslant 1\right\}
$$

and $M=\sup \left\{\|x\|_{\infty} ; x \in X,\|x\| \leqslant 1\right\}$ which we assume to be finite. Then for some absolute constant $C$ and all $0<\varepsilon<1 / 2$,

$$
\begin{equation*}
N(X, \varepsilon) \leqslant C \varepsilon^{-2}\left(\left(\int_{1}^{M}\left(\frac{\log E\left(B_{X}, B_{\infty}, t\right)}{t}\right)^{1 / 2} d t\right)^{2}+n \log \varepsilon^{-1}\right) \tag{5.1}
\end{equation*}
$$

Proof. Let $l=\left[\log _{2} M\right]+1$ and for $s=0,1, \ldots, l$ let $\mathscr{A}_{s}$ be a subset of $B_{X}$ so that

$$
\begin{equation*}
B_{X} \subset \bigcup_{y \in \mathscr{A}_{s}}\left\{y+2^{s} B_{\infty}\right\}, \quad \overline{\mathscr{A}}_{s} \leqslant \tilde{E}\left(B_{X}, B_{\infty}, 2^{s}\right) \tag{5.2}
\end{equation*}
$$

For $s=l$ we choose $\mathscr{A}_{l}$ to consist just of the origin 0 . Let $\mathscr{F}$ be an $\varepsilon$-net in $B_{X}$ so that

$$
\begin{equation*}
\log \underset{\mathscr{F}}{\mathscr{F}} \leqslant c_{1} n \log \varepsilon^{-1} \tag{5.3}
\end{equation*}
$$

(cf. Lemma 2.4). For every $x \in \mathscr{F}$ and $0 \leqslant s \leqslant l$ let $y_{s, x} \in \mathscr{A}_{s}$ satisfy $\left\|x-y_{s, x}\right\|_{\infty} \leqslant 2^{s}$. We clearly have

$$
\begin{equation*}
|x(\omega)|=\left(|x(\omega)|-\left|y_{0, x}(\omega)\right|\right)+\sum_{s=1}^{l}\left(\left|y_{s, x}(\omega)\right|-\left|y_{s-1, x}(\omega)\right|\right), \quad x \in X, \quad \omega \in \Omega \tag{5.4}
\end{equation*}
$$

We are going to apply the empirical method to the following collection of subsets of $L_{\infty}(\Omega, \mu)$

$$
\mathscr{C}_{0}=\left\{|x|-\left|y_{0, x}\right| ; x \in \mathscr{F}\right\}, \quad \mathscr{C}_{s}=\left\{\left|y_{s, x}\right|-\left|y_{s-1, x}\right| ; x \in \mathscr{F}\right\}, \quad 1 \leqslant s \leqslant l .
$$

Since

$$
\left\|\left|y_{s, x}\right|-\left|y_{s-1, x}\right|\right\|_{\infty} \leqslant\left\|y_{s, x}-x\right\|_{\infty}+\left\|y_{s-1, x}-x\right\|_{\infty} \leqslant 2^{s}+2^{s-1} \leqslant 2^{s+1}
$$

we have

$$
\begin{equation*}
\|f\|_{1} \leqslant 2, \quad\|f\|_{\infty} \leqslant 2^{s+1}, \quad f \in \mathscr{C}_{s}, \quad 0 \leqslant s \leqslant l \tag{5.5}
\end{equation*}
$$

Also by (4.1), (5.2) and (5.3)

$$
\begin{gather*}
\log \overline{\mathscr{E}}_{0} \leqslant c_{1} n \log \varepsilon^{-1},  \tag{5.6}\\
\log \overline{\overline{\mathscr{E}}}_{s} \leqslant \log \overline{\bar{A}}_{s}+\log \overline{\mathscr{A}}_{s-1} \leqslant 2 \log E\left(B_{X}, B_{\infty}, 2^{s-2}\right), \quad 1 \leqslant s \leqslant l .
\end{gather*}
$$

Our aim it to find $\left\{\omega_{j}\right\}_{j=1}^{N}$ in $\Omega$ where $N$ is given by the right hand side of (5.1) so that

$$
\begin{gather*}
\left|N^{-1} \sum_{j=1}^{N} f\left(\omega_{j}\right)-\int_{\Omega} f(\omega) d \mu(\omega)\right| \leqslant \varepsilon, \quad f \in \mathscr{E}_{0}  \tag{5.7}\\
\left|N^{-1} \sum_{j=1}^{N} f\left(\omega_{j}\right)-\int_{\Omega} f(\omega) d \mu(\omega)\right| \leqslant \varepsilon_{s}, \quad f \in \mathscr{C}_{s}, \quad 1 \leqslant s \leqslant l \text { with } \sum_{s=1}^{l} \varepsilon_{s} \leqslant \varepsilon \tag{5.8}
\end{gather*}
$$

Once (5.7) and (5.8) hold we get by (5.4) that

$$
\left|N^{-1} \sum_{j=1}^{N}\right| x\left(\omega_{j}\right)|-\|x\|| \leqslant 2 \varepsilon, \quad x \in \mathscr{F}
$$

and an application of Lemma 2.5 will conclude the proof.
To ensure that (5.7) and (5.8) hold it is enough to verify that the measure of the set of $\left\{\omega_{j}\right\}_{j=1}^{N}$ in $\Omega^{N}$ for which at least one of the inequalities fails is less than 1 . In view of Lemma 2.1 this will be the case in particular if

$$
\begin{gather*}
2 \overline{\mathscr{E}}_{0} \exp \left(-\varepsilon^{2} N / 32\right)<1 / 2  \tag{5.9}\\
2 \sum_{s=1}^{i} \overline{\mathscr{E}}_{s} \exp \left(-\varepsilon_{s}^{2} N / 32 \cdot 2^{s+1}\right)<1 / 2 \tag{5.10}
\end{gather*}
$$

(We apply Lemma 2.1 to the functions $f-\int f d \mu, f \in \mathscr{C}_{s}$; these functions have an $L_{1}$ norm bounded by 4 and an $L_{\infty}$ norm bounded by $2^{s+1}$.) By (5.6) it follows that (5.9) is ensured once

$$
\begin{equation*}
N \geqslant c_{2} n \varepsilon^{-2} \log \varepsilon^{-1} \tag{5.11}
\end{equation*}
$$

Note that since obviously $E\left(B_{X}, B_{\infty}, \lambda\right) \geqslant M / \lambda$, it follows from (5.6) that (5.10) is ensured once we have for a suitable $c_{3}>0$

$$
\exp \left(-\varepsilon_{s}^{2} N / 32 \cdot 2^{s+1}\right) \geqslant c_{3} E\left(B_{X}, B_{\infty}, 2^{s-2}\right)^{4}, \quad 1 \leqslant s \leqslant l
$$

i.e.,

$$
\varepsilon_{s} N^{1 / 2} \geqslant c_{4}\left(\log E\left(B_{X}, B_{\infty}, 2^{s-2}\right)\right)^{1 / 2} 2^{s / 2}, \quad 1 \leqslant s \leqslant l
$$

Since we also need to have that $\Sigma_{s=1}^{\prime} \varepsilon_{s} \leqslant \varepsilon$, and $E\left(B_{X}, B_{\infty}, t\right)$ is a monotone decreasing function of $t$, we get that what will ensure (5.10) is

$$
\begin{equation*}
\varepsilon N^{1 / 2} \geqslant c_{4} \sum_{s=1}^{l} \log E\left(B_{X}, B_{\infty}, 2^{s-2}\right) 2^{s / 2} \geqslant c_{5} \int_{1}^{M}\left(\frac{\log E\left(B_{X}, B_{\infty}, t\right)}{t}\right)^{1 / 2} d t \tag{5.12}
\end{equation*}
$$

Since the right hand side of (5.1) verifies both (5.11) and (5.12) the proposition is proved.

We are now ready to prove Theorem 1.2 in the following, more precise, formulation.

Theorem 5.2. Let $X$ be an dimensional subspace of $L_{1}(0,1)$. Then for some absolute constant $c$

$$
\begin{equation*}
N(X, \varepsilon) \leqslant c \varepsilon^{-2} \log \left(n \varepsilon^{-1}\right)(\log n)^{2} n, \quad \varepsilon>0 \tag{5.13}
\end{equation*}
$$

Proof. By Lemma 2.3 and Corollary 4.8 there is for every $\varepsilon>0$ a probability space $(\Omega, v)$ and a subspace $\tilde{X}$ of $L_{1}(\Omega, v)$ so that $d(X, \tilde{X}) \leqslant 1+\varepsilon, \sup \left\{\|x\|_{\infty} ; x \in \tilde{X},\|x\|=1\right\} \leqslant 2 n$ and

$$
\log E\left(B_{\bar{X}}, B_{\infty}, t\right) \leqslant c_{1}\left[\log n \log \left(n \varepsilon^{-1}\right)\right]^{1 / 2} n / t, \quad t>1
$$

Hence by (5.1)

$$
N(\tilde{X}, \varepsilon) \leqslant c_{2} \varepsilon^{-2} n\left(\left(\int_{1}^{2 n} \frac{d t}{t}\right)^{2} \cdot\left[\log n \cdot \log \left(n \varepsilon^{-1}\right)\right]^{1 / 2}+\log \varepsilon^{-1}\right)
$$

and this implies (5.13).

We pass now to estimates on $N(X, \varepsilon)$ for spaces for which we have information on $T_{p}(X)$ for some $p>1$. Our aim is to prove Theorem 1.1. We shall combine in the proof
the weak $L_{p}$ estimates given by Proposition 3.1 with the entropy estimates from section 4. The strategy is just to combine the proof of Lemma 3.3. with that of Proposition 5.1. While the proof involves no additional ideas to those already used above the task of keeping track of all the parameters leads to some quite lengthy (but elementary) computations.

Lemma 5.3. Let $X$ be an $n$-dimensional subspace of $l_{1}^{N}$. Let $1<p \leqslant 2,0<\varepsilon<1 / 2$ and assume that $N \geqslant n \varepsilon^{-2}$. Then there is a subspace $Y$ of $l_{1}^{\dot{N}}$ with $d(X, Y) \leqslant 1+\varepsilon$ provided that

$$
\begin{equation*}
\tilde{N} \geqslant c n \varepsilon^{-2}(\log T)^{1 / 2}(p-1)^{-3}\left(\log (N T / n(p-1))^{5 / 2}\right. \tag{5.14}
\end{equation*}
$$

where $c$ is a universal constant and $T=T_{p}(X)+1$.

Proof. By Lemmas 2.3 and 3.2 we may assume that there is a probability measure $\mu$ on $\{1,2, \ldots, N\}$ so that $\mu\{i\}>1 / 2 N$ for every $i$ and

$$
\begin{equation*}
\|x\|_{L_{\infty}(\mu)} \leqslant c_{2} n\|x\|, \quad\|x\|_{L_{p, \infty}(\mu)} \leqslant c_{2} T_{p}(X)(p-1)^{-1}\|x\|, . x \in X . \tag{5.15}
\end{equation*}
$$

In addition, in view of Lemma 4.3 and Proposition 4.6, we may also assume that for every $q>1$

$$
\begin{equation*}
\log E\left(B_{X}, B_{q}, t\right) \leqslant c_{3} K(X) n q^{1 / 2} / t \leqslant c_{4}\left(\frac{\log \left(1+T_{p}(X)\right)}{p-1}\right)^{1 / 2} q^{1 / 2} n / t \tag{5.16}
\end{equation*}
$$

Note that, by (5.15), $\|x\|_{L_{q}(\mu)} \leqslant c_{2} n^{1-1 / q}\|x\|$ for every $x \in X$ and thus for $t>c_{2} n^{1-1 / q}$, $E\left(B_{X}, B_{q}, t\right)=1$. Let $\mathscr{F}$ be an $\varepsilon$-net in the unit ball of $X$ satisfying (5.3). We fix now a $q \geqslant 2$ and put

$$
\begin{equation*}
\alpha=8 c_{2} T_{p}(X) /(p-1), \quad \beta=8 c_{4}\left(\frac{\log \left(1+T_{p}(X)\right)}{p-1}\right)^{1 / 2} \cdot q^{1 / 2}, \quad l=\left[\log c_{2} n^{1-1 / q}\right] \tag{5.17}
\end{equation*}
$$

For $0 \leqslant s \leqslant l$ we choose now a set $G_{s}$, in $B_{X}$ so that

$$
\begin{equation*}
B_{X} \subset \bigcup_{y \in G_{s}}\left\{y+2^{s} B_{q}\right\}, \quad \overline{\bar{G}}_{s} \leqslant \tilde{E}\left(B_{X}, B_{q}, 2^{s}\right) \tag{5.18}
\end{equation*}
$$

The set $G_{l}$ consists just of the origin. For $x \in \mathscr{F}, 0 \leqslant s \leqslant l$ let $y_{s, x} \in G_{s}$ satisfy $\left\|x-y_{s, x}\right\|_{q} \leqslant 2^{s}$. Put next

$$
\mathscr{C}_{0}=\left\{|x|-\left|y_{0, x}\right| ; x \in \mathscr{F}\right\}, \quad \mathscr{C}_{s}=\left\{\left|y_{s, x}\right|-\left|y_{s-1, x}\right| ; x \in \mathscr{F}\right\}, \quad 1 \leqslant s \leqslant l .
$$

By (4.1), (5.16), (5.17) and (5.18) we get that

$$
\begin{equation*}
\log \overline{\mathscr{E}}_{s} \leqslant \beta n 2^{-s}, \quad 1 \leqslant s \leqslant l . \tag{5.19}
\end{equation*}
$$

By (5.15) we get that for every $z \in \mathscr{E}_{s}, 0 \leqslant s \leqslant l$

$$
\begin{equation*}
\|z\|_{L_{p, \infty}(\mu)} \leqslant 2, \quad\|z\|_{L_{q}(\mu)} \leqslant 2^{s+1}, \quad\|z\|_{L_{\infty}(\mu)} \leqslant \min \left((2 N)^{1 / q} 2^{s+1}, 2 c_{2} n\right) \tag{5.20}
\end{equation*}
$$

As in the proof of Proposition 5.1 we shall apply the empirical method to all the sets $\mathscr{E}_{s}, 0 \leqslant s \leqslant l$. We find below positive numbers $\left\{\theta_{s}\right\}_{s=1}^{l}$ and $\left\{\delta_{s}\right\}_{s=1}^{l}$ and certain expressions $N_{1, q}, N_{2, q}$ and $N_{3, q}$ so that

$$
\begin{equation*}
\sum_{s=1}^{l} \theta_{s} \leqslant \varepsilon, \quad \sum_{s=1}^{l} \delta_{s}<1 / 2 \tag{5.21}
\end{equation*}
$$

and so that for $\tilde{N} \geqslant \max \left(N_{1, q}, N_{2, q}, N_{3, q}\right)$ we have

$$
\begin{gather*}
\operatorname{Prob}\left\{\left|\tilde{N}^{-1} \sum_{j=1}^{\tilde{N}} z\left(i_{j}\right)-\int z d \mu\right|>\varepsilon \text { for at least one } z \in \mathscr{E}_{0}\right\}<1 / 2  \tag{5.22}\\
\operatorname{Prob}\left\{\left|\tilde{N}^{-1} \sum_{j=1}^{\bar{N}} z\left(i_{j}\right)-\int z d \mu\right|>\theta_{s} \text { for at least one } z \in \mathscr{C}_{s}\right\}<\delta_{s}, \quad 1 \leqslant s \leqslant l, \tag{5.23}
\end{gather*}
$$

where the probability of $\left\{i_{j}\right\}_{j=1}^{\mathcal{N}}$ is taken of course with respect to $\mu^{N}$. Once we verify this the proof will be concluded by showing that the right hand side of (5.14) exceeds $\max \left(N_{1, q}, N_{2, q}, N_{3, q}\right)$ for a suitably chosen $q$.

We start with (5.22). As in the proof of Lemma 3.3 we cut each $z$ in $\mathscr{E}_{0}$ into its level sets $A_{k}(z)=\left\{i ; 2^{k-1}<|z(i)| \leqslant 2^{k}\right\}$ and consider instead of $\mathscr{E}_{0}$ the $\sim \log n$ sets $\mathscr{E}_{0, k}$ consisting of $z \cdot \chi_{A_{k}(z)}$. We apply the empirical method to all those sets using the reasoning of Lemma 3.3 but with $L_{q}$ estimates (instead of $L_{p, \infty}$ estimated used there). Thus, for example, (3.16) becomes in our setting (note that $\overline{\mathscr{F}}=\overline{\mathscr{F}}_{0}$ )

$$
2^{k_{0}} \leqslant c_{5}\left(\frac{N}{\log \mathscr{F}}\right)^{1 / q}\left(\log \left(\frac{N}{\log \mathscr{F}}\right)\right)^{1 / q} .
$$

We get from this argument that (5.22) holds provided that

$$
\begin{equation*}
\tilde{N} \geqslant N_{1, q}=c_{6} n \varepsilon^{-2} \log \frac{1}{\varepsilon}\left(\frac{N}{n} \log \frac{N}{n}\right)^{1 / q} \tag{5.24}
\end{equation*}
$$

We pass now to sets $\mathscr{E}_{s}$ for $1 \leqslant s \leqslant l$. Again we define $\mathscr{E}_{s, k}=\left\{z \chi_{A_{k}(z)} ; z \in \mathscr{E}_{s}\right\}$ and apply the empirical method to these sets. By using the $L_{q}$ and $L_{\infty}$ estimates from (5.20) the proof of Lemma 3.3 shows that (5.23) is satisfied for a given $s$ if

$$
\begin{equation*}
\tilde{N} \geqslant c_{7} \theta_{s}^{-2} 2^{s} \log \overline{\overline{\mathscr{E}}}_{s}\left(\frac{N}{\log \overline{\overline{\mathscr{C}}}_{s}} \cdot \log \frac{N}{\log \overline{\mathscr{E}}_{s}}\right)^{1 / q} \tag{5.25}
\end{equation*}
$$

and

$$
\begin{equation*}
\delta_{s}>\exp \left(-c_{8} \theta_{s}^{2} \tilde{N} / 2^{s} N^{1 / q}\right) \tag{5.26}
\end{equation*}
$$

Similarly, by using the $L_{p, \infty}$ estimates instead of the $L_{q}$ estimates the same argument gives that (5.23) holds for a given $s$ if

$$
\begin{equation*}
\tilde{N} \geqslant c_{7} \theta_{s}^{-2} \alpha(p-1)^{-2} \log \overline{\mathscr{E}}_{s}\left(\frac{N}{\log \mathscr{E}_{s}} \log \frac{N}{\log \overline{\bar{E}}_{s}}\right)^{1 / p} \tag{5.27}
\end{equation*}
$$

and

$$
\begin{equation*}
\delta_{s}>\exp \left(-c_{8} \theta_{s}^{2}(p-1)^{2} \tilde{N} / 2^{s} N^{1 / q}\right) \tag{5.28}
\end{equation*}
$$

Our choice of $\left\{\theta_{s}\right\}_{s=1}^{l}$ will be so that in particular $\theta_{s} \geqslant c_{9} \varepsilon(3 / 2)^{s-l}, 1 \leqslant s \leqslant l$, for some positive (small) absolute constant $c_{9}$. This requirement does not contradict the first inequality in (5.21) and allows us to choose the $\delta_{s}$ so that the second inequality in (5.21) holds as well as (5.28) (and thus also (5.26)) for every $1 \leqslant s \leqslant l$, provided that

$$
\begin{equation*}
\tilde{N} \geqslant N_{2, q}=c_{10} n \varepsilon^{-2}(p-1)^{-2}\left(\frac{N}{n}\right)^{1 / q} . \tag{5.29}
\end{equation*}
$$

We now let $s_{0}$ be defined by the equation

$$
\begin{equation*}
\frac{\alpha}{(p-1)^{2}}\left(\frac{N}{\log \overline{\bar{E}}_{\left[s_{0}\right]}} \log \frac{N}{\log \overline{\bar{E}}_{\left[s_{0}\right]}}\right)^{1 / p}=2^{s_{0}}\left(\frac{N}{\log \overline{\bar{E}}_{\left[s_{0}\right]}} \log \frac{N}{\log \overline{\tilde{E}}_{\left[s_{0}\right]}}\right)^{1 / q} \tag{5.30}
\end{equation*}
$$

and define $\theta_{s}$, so that for $h>0$

$$
\begin{equation*}
\theta_{\left[s_{0}\right]+h} \geqslant \theta_{\left[s_{0}\right]}\left(\frac{3}{2}\right)^{-\frac{h}{2}(1-1 / p)}, \quad \theta_{\left[s_{0}\right]^{-h}} \geqslant \theta_{\left[s_{0}\right]}\left(\frac{3}{2}\right)^{-h / 2 q} \tag{5.31}
\end{equation*}
$$

Such a choice can be made (consistent with our previous requirement on $\theta_{s}$ ) so that $\Sigma_{s=1}^{l} \theta_{s} \leqslant \varepsilon$ provided that $\theta_{\left[s_{0}\right]} \geqslant c_{11} \varepsilon \min \left((p-1), q^{-1}\right)$. It is clear from (5.30) and (5.31) that if $\tilde{N}$ is such that (5.27) holds for $\left[s_{0}\right]$ then (5.27) holds for all $s \geqslant\left[s_{0}\right]$ and (5.25) holds for $s<\left[s_{0}\right]$.

In other words in order to ensure that for every $1 \leqslant s \leqslant l$ either (5.25) or (5.27) holds for our choice of $\left\{\theta_{s}\right\}_{s=1}^{l}$ it suffices that

$$
\begin{equation*}
\tilde{N} \geqslant c_{12} \varepsilon^{-2} \alpha(p-1)^{-2}\left(\max \left((p-1)^{-1}, q\right)\right)^{2} \log \overline{\bar{E}}_{\left[s_{0}\right]}\left(\frac{N}{\log \overline{\bar{E}}_{\left[s_{0}\right]}} \cdot \log \frac{N}{\bar{E}_{\left[s_{0}\right]}}\right)^{1 / p} \tag{5.32}
\end{equation*}
$$

A direct computation using (5.17), (5.19) and (5.30) shows that (5.32) holds once

$$
\begin{equation*}
\tilde{N} \geqslant N_{3, q}=c_{13} n \varepsilon^{-2} \max \left(\frac{1}{(p-1)^{2}}, q^{2}\right)\left(\frac{q \log T}{p-1}\right)^{1 / 2}\left[\left(\log \frac{N T}{n(p-1)}\right)\left(\frac{N T}{n(p-1)^{4}}\right)\right]^{r} \tag{5.33}
\end{equation*}
$$

where $T=T_{p}(X)+1$ and $r=p /(p q-q+p)$.
For $q=(p-1)^{-1} \log (N T / n(p-1))$ the right hand side of (5.14) satisfies all three conditions on $\tilde{N}$ namely (5.24), (5.29) and (5.33).

As in section 3 we can now use iteration and get a good estimate for $N(X, \varepsilon)$ for every $X \subset L_{1}(0,1)$. In particular we get the following refined version of Theorem 1.1.

Theorem 5.4. Let $X$ be an n-dimensional subspace of $L_{1}(0,1)$, let $0<\varepsilon<1 / 2$, let $1<p \leqslant 2$ and $\tau>0$. Then

$$
\begin{equation*}
N(X, \varepsilon) \leqslant c(\tau) n \varepsilon^{-2}\left(\log \left(T_{p}(X)+1\right)\right)^{1 / 2}(p-1)^{-3-\tau}\left(\log \left(T_{p}(X) / \varepsilon(p-1)\right)\right)^{5 / 2+\tau} \tag{5.34}
\end{equation*}
$$

where $c(\tau)$ is a function of $\tau$ only.

Proof. Put $\varphi(\varepsilon)=N(X, \varepsilon) / n, \alpha=\left(\log \left(T_{p}(X)+1\right)\right)^{1 / 2}(p-1)^{-3}, \beta=\log \left(\left(T_{p}(X)+1\right) /(p-1)\right)$. By Lemma 5.3 we have for every integer $k$ that whenever $\varphi\left(\varepsilon 2^{-k}\right) \geqslant 2^{2 k} \varepsilon^{-2}$ then

$$
\begin{equation*}
\varphi\left(\varepsilon 2^{-(k-1)}\right) \leqslant c 2^{2 k} \varepsilon^{-2} \alpha\left(\beta+\log \varphi\left(\varepsilon 2^{-k}\right)\right)^{5 / 2} \tag{5.35}
\end{equation*}
$$

By a routine calculation (5.34) follows from (5.35).

Remark. If we take $p-1=(\log n)^{-1}$ in Theorem 5.4 we get a slightly weaker version of Theorem 5.2. Indeed for every $n$-dimensional subspace $X$ of $L_{1}(0,1), T_{p}(X) \leqslant c_{1}$ (for $p=1+(\log n)^{-1}$ ). Hence the dependence of $N(X, \varepsilon)$ on $n$ we obtain is $n(\log n)^{3+\tau}$ (we get $n(\log n)^{3}$ in Theorem 5.2).

## 6. The approximation of Euclidean balls

In the previous sections we treated the question of embedding arbitrary subspaces of $L_{1}(0,1)$ in $l_{1}^{N}$. In [F.L.M.] the question of embedding Euclidean spaces into arbitrary Banach spaces was considered. This section is devoted to some results concerning the case which is common to both of these studies, i.e., the embedding of $l_{2}^{n}$ in $t_{1}^{N}$, or in geometric language, the approximation of Euclidean balls by zonotopes. While this case has been examined in several papers by now, and is certainly simpler than the general problems treated in the previous sections or in [F.L.M.], there are still (as we shall see) some open problems related to it.

In [Be.Mc.] the authors raise the question of how well $B^{n}$ can be approximated by a sum of segments of equal length. The empirical distribution method can be easily applied also for treating this question. We start by proving a simple variant of Lemma 2.1. This variant involves Orlicz spaces so we recall the definition of such spaces. Let $(\Omega, \mu)$ be a probability space and let $\Phi(t)$ be a convex increasing function on $[0, \infty)$ so that $\Phi(0)=0$ and $\lim _{t \rightarrow \infty} \Phi(t)=\infty$. We denote by $L_{\Phi}(\mu)$ the space of all measurable realvalued functions $f$ on $\Omega$ so that $\int_{\Omega} \Phi(|f| / \lambda) d \mu<\infty$ for some $\lambda>0$ and put

$$
\begin{equation*}
\|f\|_{L_{\Phi}(\mu)}=\inf \left\{\lambda>0 ; \int_{\Omega} \Phi(|f| / \lambda) d \mu \leqslant 1\right\} \tag{6.1}
\end{equation*}
$$

We shall be concerned here only with the two functions

$$
\begin{equation*}
\psi_{1}(t)=e^{t}-1, \quad \psi_{2}(t)=e^{t^{2}}-1, \quad t \geqslant 0 \tag{6.2}
\end{equation*}
$$

(besides of course the functions $t^{p}$ which give rise to the $L_{p}(\mu)$ spaces).

Lemma 6.1. Let $\left\{g_{j}\right\}_{j=1}^{N}$ be independent random variables with mean 0 on some probability space $(\Omega, \mu)$. Assume that $\left\|g_{j}\right\|_{L_{\psi_{1}}(\mu)} \leqslant M$ for some constant $M$ and every $1 \leqslant j \leqslant N$. Then for $0<\varepsilon<4 M$

$$
\begin{equation*}
\operatorname{Prob}\left\{\left|\sum_{j=1}^{N} g_{j}\right| \geqslant \varepsilon N\right\} \leqslant 2 \exp \left(-\varepsilon^{2} N / 16 M^{2}\right) \tag{6.3}
\end{equation*}
$$

Proof. We may clearly assume that $M=1$. For every $g$ with $\left\|g_{j}\right\|_{L_{\psi_{1}}(\mu)} \leqslant 1$ we have $\int_{\Omega} \exp (|g|) d \mu \leqslant 2$ and hence

$$
\int_{\Omega}|g|^{k} d \mu \leqslant k!\int_{\Omega} \exp (|g|) d \mu \leqslant 2 \cdot k!, \quad 1 \leqslant k<\infty
$$

Consequently, if in addition $\int_{\Omega} g d \mu=0$, we get for $0<\lambda \leqslant 1 / 2$

$$
\int_{\Omega} \exp (\lambda g) d \mu \leqslant 1+\sum_{k=2}^{\infty} \int_{\Omega}|\lambda g|^{k} d \mu / k!\leqslant 1+2 \sum_{k=2}^{\infty} \lambda^{k} \leqslant 1+4 \lambda^{2} .
$$

Hence

$$
\begin{aligned}
\operatorname{Prob}\left\{\sum_{j=1}^{N} g_{j}>\varepsilon N\right\} \exp (\varepsilon \lambda N) & \leqslant \int \exp \left(\lambda \sum_{j=1}^{N} g_{j}\right) d \mu=\prod_{j=1}^{N} \int \exp \left(\lambda g_{j}\right) d \mu \\
& \leqslant\left(1+4 \lambda^{2}\right)^{N} \leqslant \exp \left(4 \lambda^{2} N\right) .
\end{aligned}
$$

By taking $\lambda=\varepsilon / 8(\leqslant 1 / 2)$ we get

$$
\operatorname{Prob}\left\{\sum_{j=1}^{N} g_{j}>\varepsilon N\right\} \leqslant \exp \left(-\varepsilon^{2} N / 16\right)
$$

and this implies (6.3).

Remark. The lemma obviously holds also if we take $M=\max \left\|g_{j}\right\|_{\psi_{\psi_{2}}(\mu)}$. In this form we shall actually apply it below. There are however instances which are of some geometric interest in which an $L_{\psi_{1}}$ estimate on the functions involved is available while no $L_{\psi_{2}}$ estimate holds. We shall discuss an example of such a situation in [B.L.M.2].

Proposition 6.2. Let $B^{n}$ be the Euclidean ball and let $0<\varepsilon<1 / 2$. Then for

$$
\begin{equation*}
N \geqslant c n \varepsilon^{-2} \log \varepsilon^{-1} \tag{6.4}
\end{equation*}
$$

there exist $N$ segments $\left\{I_{j}\right\}_{j=1}^{N}$ of equal length so that

$$
\begin{equation*}
(1-\varepsilon) B^{n} \subset \sum_{j=1}^{N} I_{j} \subset(1+\varepsilon) B^{n} . \tag{6.5}
\end{equation*}
$$

Proof. Consider $S^{n-1}=\partial\left(B^{n}\right)$ with the normalized rotation invariant measure $\sigma_{n}$. Let $X^{n}$ be the subspace of $L_{1}\left(\sigma_{n}\right)$ consisting of the (restrictions to $S^{n-1}$ of the) linear functions. (We identify the elements of $X^{n}$ with the points in $\mathbf{R}^{n}$.) Put

$$
\begin{equation*}
\beta_{n}=\int_{s^{n-1}}|\langle x, u\rangle| d \sigma_{n}(u), \quad x \in S^{n-1} . \tag{6.6}
\end{equation*}
$$

Then $\beta_{n} \sim n^{-1 / 2}$. A simple computation also shows that for $x \in S^{n-1}$

$$
\int \exp \left(n\langle x, u\rangle^{2} / 4\right) d \sigma_{n}(u)=\int_{0}^{\pi / 2} \cos ^{n-2} \theta \exp \left(n \sin ^{2} \theta / 4\right) d \theta / \int_{0}^{\pi / 2} \cos ^{n-2} \theta d \theta \leqslant c_{1}
$$

for some absolute constant $c_{1}$. Hence

$$
\begin{equation*}
\left\|\beta_{n}^{-1}\langle x, u\rangle\right\|_{L_{\psi_{2}}\left(\sigma_{n}\right)} \leqslant c_{2}, \quad x \in S^{n-1} \tag{6.7}
\end{equation*}
$$

We now apply Lemma 6.1 to functions of the form

$$
\beta_{n}^{-1}\left|\left\langle x, u_{j}\right\rangle\right|-1, \quad x, u_{j} \in S^{n-1}
$$

and get by the reasoning of section 2 in view of (6.2), (6.3) and (6.7) that for a suitable choice of $\left\{u_{j}\right\}_{j=1}^{N}$ in $S^{n-1}$

$$
\begin{equation*}
\left|\left(N \beta_{n}\right)^{-1} \sum_{j=1}^{N}\right|\left\langle x, u_{j}\right\rangle|-1|<\varepsilon, \quad x \in S^{n-1} \tag{6.8}
\end{equation*}
$$

It follows from the Hahn Banach theorem (or more precisely the separation theorem) that (6.8) implies (6.5) if we take $I_{j}=\left(N \beta_{n}\right)^{-1}\left[-u_{j}, u_{j}\right]$. All these intervals have a common length $2 / N \beta_{n} . K(X)$ and (4.32) follows.

The proof of Proposition 6.2 actually enables us to obtain a more general geometric result. Proposition 6.2 asserts that we can approximate $B^{n}$ by a sum of $N=N(n, \varepsilon)$ sets each obtained from a fixed segment by a suitable rotation. It turns out that it is possible to replace the segment above by an arbitrary compact convex set.

Theorem 6.3. Let $K$ be a compact convex set in $\mathbf{R}^{n}$ and let $0<\varepsilon<1 / 2$. There are a constant $r=r(K)$ and orthogonal transformations $\left\{U_{j}\right\}_{j=1}^{N}$ on $\mathbf{R}^{n}$, with $N \leqslant c n \varepsilon^{-2} \log \varepsilon^{-1}$ so that

$$
\begin{equation*}
(1-\varepsilon) r B^{n} \subset \frac{1}{N} \sum_{j=1}^{N} U_{j} K \subset(1+\varepsilon) r B^{n} \tag{6.9}
\end{equation*}
$$

Proof. There is no loss of generality to assume that $K$ is symmetric with respect to the origin (otherwise replace $K$ by $K-K$ ). There is also no loss of generality to assume that $K$ has the origin as an interior point. Let \|| \| be the norm on $\mathbf{R}^{n}$ whose unit ball is the polar of $K$, i.e.,

$$
\|x\|=\sup \{|\langle x, u\rangle|, u \in K\}, \quad x \in \mathbf{R}^{n}
$$

We shall now proceed as in the proof of Proposition 6.2. The only difference is that the $L_{\psi_{2}}$ estimate we need now lies much deeper. This estimate is however available in the literature. Let $O^{n}$ be the group of orthogonal transformations on $\mathbf{R}^{n}$ and let $\mu_{n}$ be the normalized Haar measure on $O^{n}$. For every $x \in \mathbf{R}^{n}$ we consider $\|U x\|$ as a function on $O^{n}$. A result of Marcus and Pisier ([Ma.P.] section V2), which is based on the so called Landau-Shepp-Fernique theorem, states that there is a universal constant $c_{1}$ so that

$$
\begin{equation*}
\|U x\|_{L_{\psi_{2}}\left(\mu_{n}\right)} \leqslant c_{1}\|U x\|_{L_{1}\left(\mu_{n}\right)}, \quad x \in \mathbf{R}^{n} \tag{6.10}
\end{equation*}
$$

Observe that

$$
\begin{equation*}
\|\mid x\|\left\|=\int_{O^{n}}\right\| U x \| d \mu_{n}(U) \tag{6.11}
\end{equation*}
$$

is a rotation invariant norm on $\mathbf{R}^{n}$ and thus is a suitable multiple of the usual Euclidean norm. By applying the empirical distribution method to the functions $\|U x\|-\|x\| \mid$ on $O^{n}$ for all $x$ satisfying $|\|x \mid\|=1$ (to be precise for an $\varepsilon$ net with respect to $\|\|\cdot\||\mid$ of this set) we get as in section 2 (and in view of Lemma 6.1 and (6.10)) that for $N$ as in the statement of the theorem and for suitable $\left\{U_{j}\right\}_{j=1}^{N} \subset O^{n}$

$$
\begin{equation*}
\left|\frac{1}{N} \sum_{j=1}^{N}\left\|U_{j} x\right\|-|\|x\|| \leqslant \varepsilon\|x \mid\|, \quad x \in \mathbf{R}^{n}\right. \tag{6.12}
\end{equation*}
$$

Assertion (6.9) follows from (6.12) by a straightforward duality argument.

Remark. For $K$ a zonotope, Theorem 6.3 is an easy consequence of Proposition 6.2 and its proof. Thus in this case we do not have to use the result of Marcus and Pisier.

A more detailed study of questions related to Theorem 6.3 is presented in [B.L.M.2]. In this paper we determine in particular the number of Minkowski symmetrizations needed in order to pass from a general convex body in $\mathbf{R}^{n}$ to a body $\varepsilon$ close to a sphere.

We turn next to the second topic of this section and this is the study whether (1.4) can be improved as far as dependence on $\varepsilon$ is concerned. We have seen in the previous section that up to logarithmic factors the same estimate holds for general zonoids and again in (6.4) we got a similar estimate for the case of equal segments. These results make the question above more interesting from the geometric point of view.

Some facts concerning this question are contained in [Be.Mc.]. One of the obser-
vations in [Be.Mc] is the following. Assume that (6.5) holds with $I_{j}=\left[-y_{j}, y_{j}\right]$. Then for every choise of (n-1) indices $\left\{j_{i}\right\}_{i=1}^{n-1}$ and every choice of signs $\left\{\theta_{i}\right\}_{i=1}^{n-1}$ there is on the boundary of $\Sigma_{j=1}^{N} I_{j}$ a segment of length $2\left\|\sum_{i=1}^{n-1} \theta_{i} y_{j}\right\|_{2}$. Hence by (6.5)

$$
\begin{equation*}
(1-\varepsilon)^{2}+\left\|\sum_{i=1}^{n-1} \theta_{i} y_{j_{i}}\right\|_{2}^{2} \geqslant(1+\varepsilon)^{2} \tag{6.13}
\end{equation*}
$$

Also by (6.5) we have for every $x \in S^{n-1}$ that $\Sigma_{j=1}^{N}\left|\left\langle x, y_{j}\right\rangle\right| \geqslant 1-\varepsilon$ and hence we may choose the $\theta_{i}$ and $j_{i}, 1 \leqslant i \leqslant n$, so that

$$
\begin{equation*}
\left\|\sum_{i=1}^{n-1} \theta_{i} y_{j_{i}}\right\|_{2} \geqslant(1-\varepsilon)(n-1) / N \tag{6.14}
\end{equation*}
$$

From (6.13) and (6.14) it follows that (for $\varepsilon<1 / 4$ )

$$
\begin{equation*}
N \geqslant n / 4 \varepsilon^{1 / 2} \tag{6.15}
\end{equation*}
$$

An estimate which is better than (6.15) for a suitable range of $\varepsilon\left(\varepsilon \geqslant n^{-2}\right.$, to be precise) is given in the next proposition.

Proposition 6.4. Assume that (6.5) holds for some choice of intervals $\left\{I_{j}\right\}_{j=1}^{N}$ not necessarily of equal length. Then for some absolute constant $c>0$

$$
\begin{equation*}
N \geqslant c n^{2} /(1+n \varepsilon) \tag{6.16}
\end{equation*}
$$

Proof. We assume that $n$ is even; the case $n$ odd follows from $n$ even by projecting on a subspace with codimension 1. Consider again the numbers $\beta_{n}$ defined in (6.6). The precise value of $\beta_{n}$ is given by

$$
\begin{equation*}
\beta_{n}=\Gamma(n / 2) /(\Gamma((n+1) / 2) \Gamma(1 / 2)) \tag{6.17}
\end{equation*}
$$

Hence, by Stirling's formula

$$
\begin{equation*}
\frac{\beta_{n / 2}}{\beta_{n}}=\frac{\Gamma(n / 4) \Gamma((n+1) / 2)}{\Gamma(n / 2) \Gamma((n+2) / 4)}=\sqrt{2}\left(1+\frac{1}{4 n}+O\left(\frac{1}{n^{2}}\right)\right) \tag{6.18}
\end{equation*}
$$

We write (6.5) in the following form. For some $\left\{y_{j}\right\}_{j=1}^{N}$ in $S^{n-1}$ and some scalars $\lambda_{1} \geqslant \lambda_{2} \geqslant \ldots \geqslant \lambda_{N}>0$,

$$
\begin{equation*}
\left|\beta_{n}^{-1} \sum_{j=1}^{N} \lambda_{j}\right|\left\langle x, y_{j}\right\rangle|-1| \leqslant \varepsilon, \quad x \in S^{n-1} \tag{6.19}
\end{equation*}
$$

By integrating (6.19) with respect to $\sigma_{n}$ we get

$$
\begin{equation*}
\left|\sum_{j=1}^{N} \lambda_{j}-1\right|<\varepsilon \tag{6.20}
\end{equation*}
$$

Hence (if $\varepsilon<1 / 2$ )

$$
\begin{equation*}
\sum_{j=1}^{n / 2} \lambda_{j} \geqslant n(1-\varepsilon) / 2 N \geqslant n / 4 N \tag{6.21}
\end{equation*}
$$

Let $E \subset \mathbf{R}^{n}$ be an $n / 2$ dimensional subspace containing $\left\{y_{j}\right\}_{j=1}^{n / 2}$ and let $Q$ be the orthogonal projection on $E$.

Then for $x \in E$ with $\|x\|_{2}=1$ we get from (6.19)

$$
\left|\beta_{n}^{-1} \sum_{j=1}^{N} \lambda_{j}\right|\left\langle x, Q y_{j}\right\rangle|-1|<\varepsilon
$$

and by integrating on the unit sphere of $E$ with respect to its invariant measure it follows that

$$
\beta_{n / 2} \beta_{n}^{-1} \sum_{j=1}^{N} \lambda_{j}\left\|Q y_{j}\right\|_{2} \geqslant 1-\varepsilon
$$

By the Cauchy Schwartz inequality

$$
\begin{equation*}
\left(\sum_{j=1}^{N} \lambda_{j}\right)\left(\sum_{j=1}^{N} \lambda_{j}\left\|Q y_{j}\right\|_{2}^{2}\right) \geqslant(1-\varepsilon)^{2}\left(\beta_{n} / \beta_{n / 2}\right)^{2} \tag{6.22}
\end{equation*}
$$

Similarly, by replacing $Q$ and $E$ with $I-Q$ and $E^{\perp}$ respectively and noting that $(I-Q) y_{j}=0$ for $1 \leqslant j \leqslant n / 2$ we get

$$
\begin{equation*}
\left(\sum_{j=n / 2+1}^{N} \lambda_{j}\right)\left(\sum_{j=n / 2+1}^{N} \lambda_{j}\left\|(I-Q) y_{j}\right\|_{2}^{2}\right) \geqslant(1-\varepsilon)^{2}\left(\beta_{n} / \beta_{n / 2}\right)^{2} \tag{6.23}
\end{equation*}
$$

and in particular (if $\varepsilon<1 / 5$ and $n \geqslant 3$ )

$$
\begin{equation*}
\sum_{j=n / 2+1}^{N} \lambda_{j}\left\|(I-Q) y_{j}\right\|_{2}^{2} \geqslant 1 / 4 \tag{6.24}
\end{equation*}
$$

By (6.20)-(6.24)

$$
(1+\varepsilon)^{2} \geqslant\left(\sum_{j=1}^{N} \lambda_{j}\right)^{2}=\left(\sum_{j=1}^{N} \lambda_{j}\right)\left(\sum_{j=1}^{N} \lambda_{j}\left\|Q y_{j}\right\|_{2}^{2}+\sum_{j=1}^{N} \lambda_{j}\left\|(I-Q) y_{j}\right\|_{2}^{2}\right)
$$

$$
\begin{aligned}
& \geqslant\left(\sum_{j=1}^{n / 2} \lambda_{j}\right)\left(\sum_{j=1}^{N} \lambda_{j}\left\|(I-Q) y_{j}\right\|_{2}^{2}\right)+2(1-\varepsilon)^{2}\left(\beta_{n} / \beta_{n / 2}\right)^{2} \\
& \geqslant n / 16 N+2(1-\varepsilon)^{2}\left(\beta_{n} / \beta_{n / 2}\right)^{2} .
\end{aligned}
$$

Hence by (6.13),

$$
n / 16 N \leqslant(1+\varepsilon)^{2}-(1-\varepsilon)^{2}\left(1-1 / 2 n+O\left(n^{-2}\right)\right)=4 \varepsilon+1 / 2 n+O(\varepsilon+1 / n)^{2}
$$

and this implies (6.16).

In the notation of section 1 , (6.16) states that $N\left(B^{n}, \varepsilon\right) \geqslant c n^{2} /(1+n \varepsilon)$. We consider next the behavior of $N=N\left(B^{n}, \varepsilon\right)$ for a fixed $n$ as a function of $\varepsilon$ as $\varepsilon \rightarrow 0$. It is trivial that for $n=2$ the right value of $N$ is given by (6.15), i.e., $N \sim c \varepsilon^{-1 / 2}$. For $n=3,4$ better upper bounds for $N$ than (1.4) were obtained (using the notion of a projection body) in [Be.Mc.], namely

$$
\begin{equation*}
N\left(B^{n}, \varepsilon\right) \leqslant \gamma_{n} \varepsilon^{(1-n) / 2} \tag{6.25}
\end{equation*}
$$

We establish here a lower bound.

Theorem 6.5. For every $n \geqslant 2$ there is a positive constant $c_{n}$ so that

$$
\begin{equation*}
N\left(B^{n}, \varepsilon\right) \geqslant c_{n} \varepsilon^{2(1-n) /(n+2)} \tag{6.26}
\end{equation*}
$$

Thus for example for $n=3$ we get from (6.25) and (6.26) that

$$
c_{3} \varepsilon^{-4 / 5} \leqslant N\left(B^{3}, \varepsilon\right) \leqslant \gamma_{3} \varepsilon^{-1}
$$

We do not know the exact order of magnitude of $N\left(B^{3}, \varepsilon\right) .\left({ }^{1}\right)$
Theorem 6.5 is an immediate consequence of the following proposition.

Proposition 6.6. Let $\mu$ be a symmetric probability measure on $S^{n-1}$ supported by $N$ points. Put

$$
\begin{equation*}
h(x)=\int_{S^{n-1}}|\langle x, y\rangle| d \mu(y) . \tag{6.27}
\end{equation*}
$$

${ }^{(1)}$ See also the note "added in proof".

Then

$$
\begin{equation*}
\delta \equiv \sup _{x \in S^{n-1}} h(x)-\inf _{x \in S^{n-1}} h(x) \geqslant\left\|h-\beta_{n}\right\|_{L_{2}\left(\sigma_{n}\right)} \geqslant c(n) N^{-(n+2) / 2(n-1)} \tag{6.28}
\end{equation*}
$$

for a suitable positive constant $c(n)$.

Here $\sigma_{n}$ denotes, as before, the normalized rotation invariant measure on $S^{n-1}$ and

$$
\beta_{n}=\int_{S^{n-1}}|\langle x, y\rangle| d \sigma_{n}(x)=\int_{S^{n-1}} h(x) d \sigma_{n}(x)
$$

Proof. We shall use some results on spherical harmonics. For each $k$ let $\left\{Y_{k, j}, 1 \leqslant j \leqslant M(n, k)\right\}$ be an orthonormal basis of the spherical harmonics of degree $k$ on $S^{n-1}$. By the addition theorem (see [Mu] p.10)

$$
\begin{equation*}
\sum_{j} Y_{k, j}(\xi) Y_{k, j}(\eta)=M(n, k) P_{k}(\langle\xi, \eta\rangle), \quad \xi, \eta \in S^{n-1} \tag{6.29}
\end{equation*}
$$

where $P_{k}(t)$ is the Legendre polynomial of degree $k$ and dimension $n$, given by Rodriques' formula (see [Mu] p. 17)

$$
\begin{equation*}
P_{k}(t)=\left(-\frac{1}{2}\right)^{k} \frac{\Gamma((n-1) / 2)}{\Gamma(k+(n-1) / 2)}\left(1-t^{2}\right)^{(3-n) / 2}\left(\frac{d}{d t}\right)^{k}\left\{\left(1-t^{2}\right)^{k+(n-3) / 2}\right\} \tag{6.30}
\end{equation*}
$$

By the Funk-Hecke formula (see [Mu] p. 20)

$$
\begin{align*}
{\left[h, Y_{k, j}\right]=\int_{s^{n-1}} h(x) Y_{k, j}(x) d \sigma_{n}(x) } & =\int_{s^{n-1}}\left[\int_{s^{n-1}}|\langle x, y\rangle| Y_{k, j}(x) d \sigma_{n}(x)\right] d \mu(y) \\
& =\lambda_{k} \int_{S^{n-1}} Y_{k, j}(y) d \mu(y) \tag{6.31}
\end{align*}
$$

where $\lambda_{k}=0$ for odd $k$ and for even $k$ (by (6.30))

$$
\begin{align*}
\lambda_{k} & =\frac{\Gamma(n / 2)}{\pi^{1 / 2} \Gamma((n-1) / 2)} \int_{-1}^{1}|t|\left(1-t^{2}\right)^{(n-3) / 2} P_{k}(t) d t \\
& =\frac{2^{-k+1} \Gamma(n / 2)}{\pi^{1 / 2} \Gamma(k+(n-1) / 2)} \int_{0}^{1} t\left(\frac{d}{d t}\right)^{k}\left\{\left(1-t^{2}\right)^{k+(n-3) / 2}\right\} d t  \tag{6.32}\\
& =\frac{(-1)^{(k-2) / 2} \Gamma(n / 2) \Gamma(k-1)}{\pi^{1 / 2} 2^{k-1} \Gamma(k / 2) \Gamma((k+n+1) / 2)} \approx k^{-n / 2-1} \quad \text { for } \quad k \rightarrow \infty
\end{align*}
$$

Note that we denote the inner product in $L_{2}\left(\sigma_{n}\right)$ by $[\cdot, \cdot]$. Recall also the following identity ([Mu] p. 30, Lemma 17)

$$
\begin{equation*}
\frac{1-r^{2}}{\left(1+r^{2}-2 r t\right)^{n / 2}}=\sum_{k=0}^{\infty} M(n, k) r^{k} P_{k}(t), \quad|t| \leqslant 1, \quad 0 \leqslant r<1 . \tag{6:33}
\end{equation*}
$$

By (6.29) we deduce that

$$
\begin{equation*}
\frac{1-r^{2}}{\left(1+r^{2}-2 r\langle\xi, \eta\rangle\right)^{n / 2}}=\sum_{k=0}^{\infty} r^{k} \sum_{j} Y_{k, j}(\xi) Y_{k, j}(\eta), \quad \xi, \eta \in S^{n-1}, \quad 0 \leqslant r<1 \tag{6.34}
\end{equation*}
$$

and by integrating with respect to $\mu$ we get (using also (6.31))

$$
\begin{equation*}
\int_{S^{n-1}} \frac{1-r^{2}}{\left(1+r^{2}-2 r\langle\xi, \eta\rangle\right)^{n / 2}} d \mu(\eta)=\sum_{k=0}^{\infty} r^{k} \lambda_{k}^{-1} \sum_{j}\left[h, Y_{k, j}\right] Y_{k, j} \tag{6.35}
\end{equation*}
$$

Since

$$
h=\beta_{n}+\sum_{k=1}^{\infty} \sum_{j}\left[h, Y_{k, j}\right] Y_{k, j}
$$

we get by the definition of $\delta$ in (6.28) and Parseval's identity that

$$
\begin{equation*}
\delta^{2} \geqslant \sum_{k=1}^{\infty} \sum_{j}\left[h, Y_{k, j}\right]^{2} \tag{6.36}
\end{equation*}
$$

By our assumption $\mu=\sum_{i=1}^{N} a_{i} \delta_{\eta_{i}}$ where $\eta_{i} \in S^{n-1}$ and $\alpha_{i} \geqslant 0$ with $\Sigma_{i=1}^{N} a_{i}=1$. Hence, by (6.35) and (6.36)

$$
\begin{equation*}
\left\|1-\sum_{i=1}^{N} \frac{a_{i}\left(1-r^{2}\right)}{\left(1+r^{2}-2 r\left\langle\xi, \eta_{i}\right\rangle\right)^{n / 2}}\right\|_{L_{2}\left(\sigma_{n}\right)} \leqslant \delta \max _{k \geqslant 1}\left(r^{k} \lambda_{k}^{-1}\right) . \tag{6.37}
\end{equation*}
$$

Next observe that by (6.34) the square of the left hand side of (6.37) is equal to

$$
\sum_{i=1}^{N} \sum_{i^{\prime}=1}^{N} a_{i} a_{i^{\prime}} \frac{1-r^{4}}{\left(1+r^{4}-2 r^{2}\left\langle\eta_{i}, \eta_{i^{\prime}}\right\rangle\right)^{n / 2}}-1
$$

which is at least

$$
-1+\left(\sum_{i=1}^{N} a_{i}^{2}\right) \frac{1-r^{4}}{\left(1-r^{2}\right)^{n}} \geqslant N^{-1}(1-r)^{-n+1}-1
$$

We choose now $r$ so that

$$
\begin{equation*}
\left(1-r^{2}\right)^{-n+1}=2 N \tag{6.38}
\end{equation*}
$$

Then by (6.32) and (6.37) there is a constant $c^{1}(n)$ so that

$$
\begin{equation*}
c^{1}(n) \delta \max _{k \text { even }}\left(r^{k} k^{1+n / 2}\right) \geqslant 1 \tag{6.39}
\end{equation*}
$$

By taking $k \sim\left(1-r^{2}\right)^{-1}$ in (6.39) and using (6.38) we deduce (6.28).
It was pointed out recently by Linhart [Li] that it follows from results of J. Beck (presented in detail in [Be.C.]) on irregularities of distribution on the sphere that the upper estimate $N\left(B^{n}, \varepsilon\right) \leqslant c(n) \varepsilon^{-2}$ can be improved also for $n \geqslant 5$. The estimate given in $[\mathrm{Li}]$ is $\left.N\left(B^{n}, \varepsilon\right) \leqslant c(n) \varepsilon^{-2+2 / n}|\log \varepsilon| .{ }^{1}\right)$

## 7. Embedding into $l_{p}^{n}$

In this section we treat the question of embedding subspaces $X$ of $L_{p}(0,1)$ into $l_{p}^{N}$ with small $N$ in the case $1<p<\infty$ (the trivial case $p=2$ is excluded). As already mentioned in the introduction, the results we get show that the estimates obtained in [F.L.M.] for the case $X$ an inner product space hold in a slightly weaker form for arbitrary $X$. In general the arguments here will resemble those used in sections $2-5$ for the case $p=1$. Therefore we shall be somewhat more brief here and mainly emphasize those technical points which have to be treated differently for $p>1$.

We start with the definition of the subject of our study. For a finite-dimensional subspace $X$ of $L_{p}(0,1)$ and $\varepsilon>0$ we put

$$
\begin{equation*}
N_{p}(X, \varepsilon)=\min \left\{N ; d(X, Y) \leqslant 1+\varepsilon \text { for some } Y \subset l_{p}^{N}\right\} . \tag{7.1}
\end{equation*}
$$

The change of density lemmas we used above for $p=1$ generalize in a straightforward way to the case $p>1$. For example the lemma of Lewis (used in section 4 for $p=1$ ) reads for general $p$ as follows.

[^2]Lemma 7.1 [Le]. Let $X$ be an $n$-dimensional subspace of $L_{p}(\Omega, \mu), 1 \leqslant p<\infty$. Then there is a basis $\left\{\varphi_{i}\right\}_{i=1}^{n}$ of $X$ so that for all scalars $\left\{\lambda_{i}\right\}_{i=1}^{n}$

$$
\begin{equation*}
\int_{\Omega}\left|\sum_{i=1}^{n} \lambda_{i} \varphi_{i}\right|^{2} F^{p-2} d \mu=n^{-1}\left(\sum_{i=1}^{n} \lambda_{i}^{2}\right) \tag{7.2}
\end{equation*}
$$

where

$$
\begin{equation*}
F=\left(\sum_{i=1}^{n} \varphi_{i}^{2}\right)^{1 / 2}, \quad\|F\|_{p}=1 \tag{7.3}
\end{equation*}
$$

It follows from this lemma that if $\nu$ is the probability measure defined on $\Omega$ by $d v=F^{p} d \mu$ then the $\operatorname{map} f \rightarrow F^{-1} f$ maps $X$ into a subspace $\tilde{X}$ of $L_{p}(\Omega, v)$ which is isometric to $X$ and which has a basis $\psi_{i}=\varphi_{i} F^{-1}, 1 \leqslant i \leqslant n$ satisfying

$$
\begin{equation*}
\sum_{i=1}^{n} \psi_{i}^{2} \equiv 1, \quad\left\|\sum_{i=1}^{n} \lambda_{i} \psi_{i}\right\|_{2}^{2}=n^{-1} \sum_{i=1}^{n} \lambda_{i}^{2} \tag{7.4}
\end{equation*}
$$

An immediate consequence of (7.4) is the fact that for $f \in \tilde{X},\|f\|_{\infty} \leqslant n^{1 / 2}\|f\|_{2}$. Consequently if $1 \leqslant p<2, f \in \tilde{X}$

$$
\|f\|_{2}^{2} \leqslant\|f\|_{\infty}^{2-p}\|f\|_{p}^{p} \leqslant n^{1-p / 2}\|f\|_{2}^{2-p}\|f\|_{p}^{p}
$$

or

$$
\begin{equation*}
\|f\|_{2} \leqslant n^{(2-p) / 2 p}\|f\|_{p} \tag{7.5}
\end{equation*}
$$

and in conclusion

$$
\begin{equation*}
\|f\|_{\infty} \leqslant n^{1 / p}\|f\|_{p}, \quad 1 \leqslant p \leqslant 2, \quad\|f\|_{\infty} \leqslant n^{1 / 2}\|f\|_{p}, \quad 2 \leqslant p, \quad f \in X \tag{7.6}
\end{equation*}
$$

As pointed out by Schechtman [Sche.3] the empirical method of section 2, applied to the expressions

$$
|f(\omega)|^{p}-\int_{\Omega}|f|^{p} d v
$$

where $f$ ranges over an $\varepsilon$-net in the boundary of the unit ball of $\tilde{X}$ gives immediately, in view of (7.6), the following estimates.

$$
\begin{align*}
& N_{p}(X, \varepsilon) \leqslant c n^{2} \varepsilon^{-2} \log \varepsilon^{-1}, \quad 1 \leqslant p<2  \tag{7.7}\\
& N_{p}(X, \varepsilon) \leqslant c n^{1+p / 2} \varepsilon^{-2} \log \varepsilon^{-1}, \quad 2<p \tag{7.8}
\end{align*}
$$

Our aim is to improve on these estimates by entropy considerations or the iteration procedure.

The next proposition gives the relevant entropy estimates.

Proposition 7.2. (i) Let $X$ be an $n$-dimensional subspace of $L_{p}(0,1)$ with $2<p<\infty$ and let $\varepsilon>0$. Then there is a probability space $(\Omega, v)$ and a subspace $\tilde{X}$ of $L_{p}(\Omega, v)$ so that $d(X, \tilde{X}) \leqslant 1+\varepsilon$ and

$$
\begin{equation*}
\log E\left(\tilde{B}_{p}, \tilde{B}_{\infty}, t\right) \leqslant c \cdot p \log \left(n \varepsilon^{-1}\right) n t^{-2}, \quad 1 \leqslant t \leqslant n^{1 / 2} \tag{7.9}
\end{equation*}
$$

(ii) With $X$ as above but for $1<p<2$ there are a $v$ and $\tilde{X}$ as above so that

$$
\begin{equation*}
\log E\left(\tilde{B}_{p}, \tilde{B}_{\infty}, t\right) \leqslant c(p-1)^{-1} \log \left(n \varepsilon^{-1}\right) n t^{-p}, \quad 1 \leqslant t \leqslant n^{1 / p} \tag{7.10}
\end{equation*}
$$

Here $\tilde{B}_{r}$ denotes the unit ball of $\tilde{X}$ in the norm induced by $L_{r}(\Omega, v)$.

Proof. The proof of part (i) is by a trivial modification of arguments appearing already in section 4. If $Y$ is a subspace of $L_{p}(v)$ for some probability space, with a basis satisfying (7.4) then by the proof of Proposition 4.6 we have for every $q>p$

$$
\begin{equation*}
\log \left(\left(B_{Y}\right)_{p},\left(B_{Y}\right)_{q}, t\right) \leqslant \log E\left(\left(B_{Y}\right)_{2},\left(B_{Y}\right)_{q}, t\right) \leqslant c_{1} q n t^{-2} \tag{7.11}
\end{equation*}
$$

By Lemma 7.1 there is a probability measure $v$ on $\left\{1, \ldots, N_{p}(X, \varepsilon)\right\}$ with $v(\{i\}) \geqslant$ $1 / 2 N_{p}(X, \varepsilon)$ for every $i$ and a subspace $\tilde{X}$ of $L_{p}(v)$ so that (7.4) holds and $d(\tilde{X}, X) \leqslant$ $1+\varepsilon$. Taking $q=\log N_{p}(X, \varepsilon)$ in (7.11) and using (7.8) we get (7.9) (see the proof of Corollary 4.7).

The proof of part (ii) requires more work namely both a duality argument and interpolation. Let again $Y$ be a subspace of $L_{p}(v)$ with a basis satisfying (7.4) and denote by $B_{r}$ its unit ball in the $L_{r}(v)$ norm.

Fix $2<r$ and let $q>r$ and $0<\theta<1$ be such that

$$
\frac{1}{r}=\frac{1-\theta}{2}+\frac{\theta}{q} .
$$

Since for $f, g \in B_{2}$

$$
\|f-g\|_{r} \leqslant\|f-g\|_{2}^{1-\theta}\|f-g\|_{q}^{\theta} \leqslant 2\|f-g\|_{q}^{\theta}
$$

we get that

$$
\begin{equation*}
\log \tilde{E}\left(B_{2}, B_{r}, t\right) \leqslant \log \tilde{E}\left(B_{2}, B_{q},(t / 2)^{1 / \theta}\right) \leqslant c_{2} q n(t / 2)^{-2 r(q-2) / q(r-2)} \tag{7.12}
\end{equation*}
$$

By taking $q=r \log n$ (and remembering that only $t \leqslant n$ count) we get that

$$
\begin{equation*}
\log \tilde{E}\left(B_{2}, B_{r}, t\right) \leqslant c_{3} r n \log n\left(t / c_{3}\right)^{-2 r /(r-2)} \tag{7.13}
\end{equation*}
$$

We shall use (7.13) for estimating $E\left(B_{p}, B_{2}, t\right)$. For $k=0,1, \ldots$, let $\mathscr{E}_{k}$ be a maximal subset of $B_{p}$ with respect to the condition $\|f-g\|_{2}>8^{k} t$ for every $f \neq g$ in $\mathscr{E}_{k}\left(\right.$ if $8^{k} t>n^{(2-p) / 2 p}$ we take $\mathscr{E}_{k}=\{0\}$ ). Clearly

$$
\overline{\mathscr{E}}_{k} \geqslant \bar{E}\left(B_{p}, B_{2}, 8^{k} t\right)
$$

It is also obvious that for every $k$ there is an $h_{k} \in B_{p}$ so that if

$$
\mathscr{F}_{k}=\left\{f ; f \in \mathscr{E}_{k},\left\|f-h_{k}\right\|_{2} \leqslant 8^{k+1} t\right\}
$$

then

$$
\begin{equation*}
\mathscr{F}_{k} \geqslant \overline{\bar{G}}_{k} / \tilde{E}\left(B_{p}, B_{2}, 8^{k+1} t\right) \geqslant \tilde{E}\left(B_{p}, B_{2}, 8^{k} t\right) / \tilde{E}\left(B_{p}, B_{2}, 8^{k+1} t\right) \tag{7.14}
\end{equation*}
$$

Let $\mathscr{G}_{k}=\left\{\left(f-h_{k}\right) / 8^{k+1} t, f \in \mathscr{F}_{k}\right\}$. Then

$$
\|g\|_{2} \leqslant 1, \quad\|g\|_{p} \leqslant 2 / 8^{k+1} t, \quad g \in \mathscr{G}_{k}, \quad\left\|g-g^{\prime}\right\|_{2}>1 / 8, \quad g \neq g^{\prime} \in \mathscr{G}_{k} .
$$

Hence for $g \neq g^{\prime}$ in $\mathscr{G}_{k}$ we have, with $p^{\prime}$ denoting the conjugate exponent to $p$,

$$
8^{-2} \leqslant\left\|g-g^{\prime}\right\|_{2}^{2} \leqslant\left\|g-g^{\prime}\right\|_{p}\left\|g-g^{\prime}\right\|_{p^{\prime}} \leqslant 4\left\|g-g^{\prime}\right\|_{p^{\prime}} / 8^{k+1} t
$$

or $\left\|g-g^{\prime}\right\|_{p^{\prime}} \geqslant 2 \cdot 8^{k-2} t$. Consequently

$$
\log \tilde{E}\left(B_{2}, B_{p^{\prime}}, 8^{k-2} t\right) \geqslant \log \overline{\bar{G}}_{k}=\log \tilde{\mathscr{F}}_{k} \geqslant \log \tilde{E}\left(B_{p}, B_{2}, 8^{k} t\right)-\log \tilde{E}\left(B_{p}, B_{2}, 8^{k+1} t\right)
$$

By summing this inequality over all $k$ and using (7.13) we get that for a suitable universal $c_{4}$

$$
\begin{equation*}
\log \tilde{E}\left(B_{p}, B_{2}, t\right) \leqslant c_{4}(p-1)^{-1} n \log n \cdot\left(t / c_{4}\right)^{-2 p /(2-p)} \tag{7.15}
\end{equation*}
$$

Since for all $1<\lambda<t$

$$
\begin{equation*}
\log E\left(B_{p}, B_{q}, t\right) \leqslant \log E\left(B_{p}, B_{2}, \lambda\right)+\log E\left(B_{2}, B_{q}, t / \lambda\right) \tag{7.16}
\end{equation*}
$$

we get by (7.7), (7.11) and (7.15) and by taking a suitable $\lambda, q=\log N_{p}(X, \varepsilon)$ and $\bar{X}$ a suitable subspace of $L_{1}(v)$ where $v$ is probability measure on $\left\{1,2, \ldots, N_{p}(X, \varepsilon)\right\}$ that (7.10) holds.

Remark. If we take $q=(4 r /(r-2)) \log s$ in (7.12) (which can be done for $s>\exp ((r-2) / 4)$ ) and then follow the argument above up to (7.16) we get that for $1<p<2<q<\infty$ and say $q p /(q-p) \geqslant 1$ that

$$
\begin{equation*}
\log E\left(B_{p}, B_{q}, t\right) \leqslant c n \log t \cdot t^{-p q(q-p)}, \quad t>\exp \left(q+(p-1)^{-1}\right) \tag{7.17}
\end{equation*}
$$

This is a very precise estimate for a general finite dimensional subspace of $L_{2}(\mu)$ in a "good position" (i.e., with a basis satisfying (7.4)). In fact, Schütt computed in [Schu] the left hand side of (7.17) for the special case of $X=l_{2}^{n}$ (i.e., $\mu$ the usual probability measure on $\{1, \ldots, n\}$ and $X$ being the entire $L_{2}(\mu)$ ), and found that

$$
\log E\left(B_{p}, B_{q}, t\right) \sim n(\log t) t^{-p q /(q-p)}
$$

We are now ready to prove sharper estimates on $N_{p}(X, \varepsilon)$ than those given by (7.7) and (7.8). We start with the case $p>2$.

Theorem 7.3 Let $2<p<\infty$ and let $0<\varepsilon<1 / 2$. Then there is a constant $c(p, \varepsilon)$ so that for every $n$-dimensional subspace $X$ of $L_{p}(0,1)$

$$
\begin{equation*}
N_{p}(X, \varepsilon) \leqslant c(p, \varepsilon) n^{p / 2} \log n \tag{7.18}
\end{equation*}
$$

Proof. By Lemma 7.1 and Proposition 7.2 there is no loss of generality to assume that for some probability space ( $\Omega, v$ ), $X \subset L_{p}(\Omega, v),\|f\|_{\infty} \leqslant n^{1 / 2}\|f\|_{p}, f \in X$ and that (7.9) holds for $E\left(B_{p}, B_{\infty}, t\right)$.

Let $\mathscr{F}$ be an $\varepsilon$-net on the boundary of $B_{p}(=$ the unit ball of $X)$ with $\log \mathscr{F} \leqslant 4 n \log \varepsilon^{-1}$. For $k=1,2, \ldots, l=\left[\log n^{1 / 2} / \log (1+\varepsilon)\right]+1$ let $\mathscr{A}_{k} \subset B_{p}$ be such that

$$
B_{p} \subset \cup_{g \in \mathscr{x}_{k}}\left(g+\varepsilon(1+\varepsilon)^{k} / 3 \cdot B_{\infty}\right)
$$

and

$$
\begin{equation*}
\log \overline{\mathscr{A}}_{k} \leqslant c_{1} p \log \left(n \varepsilon^{-1}\right) n \varepsilon^{-2}(1+\varepsilon)^{-2 k} . \tag{7.19}
\end{equation*}
$$

For every $x \in \mathscr{F}$ and $1 \leqslant k \leqslant l$ let $f_{k, x} \in \mathscr{A}_{k}$ satisfy $\left\|x-f_{k, x}\right\|_{\infty} \leqslant \varepsilon(1+\varepsilon)^{k} / 3$. Put

$$
\begin{gathered}
C_{k, x}=\left\{\omega ;\left|f_{k, x}(\omega)\right| \geqslant(1+\varepsilon)^{k-1}\right\} \\
D_{k, x}=C_{k, x} \sim \bigcup_{h>k} C_{h, x}, \quad D_{0, x}=\Omega \sim \cup_{k \geqslant 1} C_{k, x}
\end{gathered}
$$

and

$$
\begin{equation*}
\tilde{x}=x \cdot \chi_{D_{0, x}}+\sum_{k=1}^{i}(1+\varepsilon)^{k} \chi_{D_{k, x}} \tag{7.20}
\end{equation*}
$$

Note that if $\omega \in C_{k, x}, k \geqslant 1$, then

$$
|x(\omega)| \geqslant(1+\varepsilon)^{k-1}-\varepsilon(1+\varepsilon)^{k} / 3>(1+\varepsilon)^{k-2}
$$

while if $\omega \notin C_{k, x}$

$$
|x(\omega)| \leqslant(1+\varepsilon)^{k-1}+\varepsilon(1+\varepsilon)^{k} / 3<(1+\varepsilon)^{k+1}
$$

Hence for every $\omega \in D_{k, x}, k \geqslant 1$,

$$
(1+\varepsilon)^{k-2} \leqslant|x(\omega)| \leqslant(1+\varepsilon)^{k+2}
$$

while for $\omega \in D_{0, x},|x(\omega)| \leqslant(1+\varepsilon)^{2}$. It follows that

$$
\begin{equation*}
(1+\varepsilon)^{-2} \leqslant|\tilde{x}(\omega)| /|x(\omega)| \leqslant(1+\varepsilon)^{2}, \quad x \in \mathscr{F}, \quad \omega \in \Omega \tag{7.21}
\end{equation*}
$$

As usual, we prove the theorem by using the empirical distribution method. We shall work with the functions

$$
|\tilde{x}(\omega)|^{p}-\int_{\Omega}|\tilde{x}(\omega)|^{p} d v, \quad x \in \mathscr{F}
$$

In view of (7.21) it follows that if we show that for suitable $\left\{\omega_{i}\right\}_{i=1}^{N}$

$$
\begin{equation*}
\left.\left|N^{-1} \sum_{i=1}^{N}\right| \tilde{x}\left(\omega_{i}\right)\right|^{p}-\int_{\Omega}|\tilde{x}(\omega)|^{p} d v \mid<\varepsilon, \quad x \in \mathscr{F} \tag{7.22}
\end{equation*}
$$

then a similar statements holds if $\tilde{x}$ is replaced by $x$ provided we replace $\varepsilon$ by $4 p \varepsilon$. In order to prove (7.22) it suffices to verify that

$$
\begin{equation*}
\left.\left|N^{-1} \sum_{i=1}^{N}\right| x\left(\omega_{i}\right)\right|^{p} \chi_{D_{0, x}}\left(\omega_{i}\right)-\int_{D_{0, x}}|x(\omega)|^{p} d v \mid<\varepsilon_{0}, \quad x \in \mathscr{F} \tag{7.23}
\end{equation*}
$$

$$
\begin{equation*}
\left|N^{-1}(1+\varepsilon)^{p k}\left(\sum_{i=1}^{N} \chi_{D_{k, x}}\left(\omega_{i}\right)-N v\left(D_{k, x}\right)\right)\right|<\varepsilon_{k}, \quad 1 \leqslant k \leqslant l, \quad x \in \mathscr{F} \tag{7.24}
\end{equation*}
$$

where

$$
\begin{equation*}
\sum_{k=0}^{l} \varepsilon_{k} \leqslant \varepsilon \tag{7.25}
\end{equation*}
$$

Let $\mathscr{B}_{k}, 1 \leqslant k \leqslant l$ be the collection of all sets of form $D_{k, x}$. It follows from the definition that

$$
\begin{equation*}
\log \overline{\mathscr{B}}_{k} \leqslant \sum_{h=k}^{l} \log \overline{\mathscr{A}}_{h} \leqslant c_{2} p \log \left(n \varepsilon^{-1}\right) n \varepsilon^{-3}(1+\varepsilon)^{-2 k} \tag{7.26}
\end{equation*}
$$

In view of Lemma 2.1, in order to ensure that (7.22) holds for some $\left\{\omega_{i}\right\}_{i=1}^{N}$, it is enough to have that

$$
\begin{equation*}
\overline{\mathscr{\mathscr { F }}} \exp \left(-\varepsilon_{0}^{2} N / 8(1+\varepsilon)^{2 p}\right)+\sum_{k=1}^{l} \overline{\mathscr{B}}_{k} \exp \left(-\varepsilon_{k}^{2} N / 8(1+\varepsilon)^{p k}\right)<\frac{1}{2} \tag{7.27}
\end{equation*}
$$

In view of (7.26) it is readily checked that (7.25) and (7.27) hold with $N=c(p, \varepsilon) n^{p / 2} \log n$ if we take $\varepsilon_{0}=c_{3} \varepsilon, \varepsilon_{l}=c_{3} \varepsilon /(p-2), \varepsilon_{k}=(1+\varepsilon)^{(2-p)(l-k) / 2} \cdot \varepsilon_{l}, 1 \leqslant k<l$.

Remark. The proof above gives that $c(p, \varepsilon)$ in (7.18) can be taken to be $c(p-2)^{-2} \varepsilon^{-5}|\log \varepsilon| \exp c p$. Some variants of the argument in the proof above give formulas for $N_{p}(X, \varepsilon)$ which improve on (7.8) as far as the dependence on $n$ is concerned without changing the dependence on $\varepsilon$. We do not go into details since for $p>1$ the dependence on $\varepsilon$ seems to be less interesting than in the case $p=1$. In the application presented in the next section only the isomorphic nature of the result (i.e., the dependence on $n$ ) is of interest.

In a completely similar manner we get, by using (7.10), the following
Theorem 7.4. Let $1<p<2$ and let $\varepsilon>0$. Then there is a constant $c(p, \varepsilon)$ so that for every $n$-dimensional subspace $X$ of $L_{p}(0,1)$

$$
\begin{equation*}
N_{p}(X, \varepsilon) \leqslant c(p, \varepsilon) n(\log n)^{3} . \tag{7.28}
\end{equation*}
$$

As in section 3 we can obtain also for $1<p<2$ a better estimate of $N_{p}(X, \varepsilon)$ (as far as dependence on $n=\operatorname{dim} X$ is concerned) if we make an extra assumption on the type of
$X$. The proof will be by iteration and the argument is identical to that used in section 3 . We just formulate here the result.

Proposition 7.5. Assume that $1<p<r \leqslant 2$ and that $X$ is an $n$-dimensional subspace of $l_{p}^{N}$. Let $0<\varepsilon<1 / 2$ and assume that $N \geqslant n \varepsilon^{-2}$. Then there is a subspace $\tilde{X}$ of $l_{p}^{\bar{N}}$ with $d(X, \tilde{X}) \leqslant 1+\varepsilon$ provided that

$$
\begin{equation*}
\tilde{N} \geqslant c n^{1-p / r} N^{p / r}\left(\log \frac{N}{n}\right)^{p / r} \varepsilon^{-2} \log \varepsilon^{-1}(r-p)^{-(p+3)} T_{r}(X)^{p} \tag{7.29}
\end{equation*}
$$

Corollary 7.6. For every $\varrho>2, \lambda \geqslant 1,1<p<r \leqslant 2$ there is a $c(p, r, \varrho, \lambda)$ so that whenever $X$ is an $n$-dimensional subspace of $L_{p}(0,1)$

$$
\begin{equation*}
N_{p}(X, \varepsilon) \leqslant c\left(p, r, \varrho, T_{r}(X)\right) n \varepsilon^{-\varrho r /(r-p)} \tag{7.30}
\end{equation*}
$$

We conclude this section with a remark concerning the role of the $l_{p}$ spaces in the discussion of this section. One may ask whether Theorems 7.3 and 7.4 or Corollary 7.6 are special cases of a general fact concerning symmetric structure. Assume that $Z$ is a Banach space with a symmetric basis $\left\{e_{i}\right\}_{i=1}^{\infty}$ and denote $Z_{n}=\operatorname{span}\left\{e_{i}\right\}_{i=1}^{n}$. For every finite dimensional subspace $X$ of $Z$ and every $\varepsilon>0$ we define $N_{Z}(X, \varepsilon)$ in an obvious analogy to (7.1). Is it true that any estimate for $N_{Z}\left(l_{2}^{n}, \varepsilon\right)$ is valid (perhaps in a slightly weaker form) for an arbitrary $n$-dimensional subspace $X$ of $Z$ ?

The answer to this question is negative. Consider for $0<\eta<1$ the space $Z(\eta)$ having a basis $\left\{e_{i}\right\}_{i=1}^{\infty}$ with

$$
\left\|\sum_{i=1}^{\infty} \lambda_{i} e_{i}\right\|=\max \left\{\left(\sum_{i=1}^{\infty} \lambda_{i}^{2}\right)^{1 / 2}, \eta \sum_{i=1}^{\infty}\left|\lambda_{i}\right|\right\}
$$

The space $Z(\eta)$ is isomorphic to $l_{1}$ and its cotype 2 constant is bounded by a constant independent of $\eta$. Hence by the results of [F.L.M.] there is a $c(\varepsilon)$ for all $\varepsilon>0$ so that $N_{Z(\eta)}\left(l_{2}^{n}, \varepsilon\right) \leqslant c(\varepsilon) n$. However $l_{1}^{2} \subset Z(\eta)$ for all $\eta>0$ and

$$
\lim _{\eta \downarrow 0} N_{Z(\eta)}\left(l_{1}^{2}, \frac{1}{10}\right)=\infty .
$$

It is also easy to construct a fixed space $Z$ of cotype 2 with a symmetric basis and subspaces $X_{n} \subset Z$ with $\operatorname{dim} X_{n}=n$ so that, e.g., $N_{Z}\left(X_{n}, 1\right) \geqslant \exp (n)$ in spite of the fact that, by [F.L.M.], $N_{Z}\left(l_{2}^{n}, \varepsilon\right) \leqslant c(\varepsilon) n$ for all $\varepsilon>0$.

## 8. Complemented subspaces

Let $X$ be an $n$-dimensional subspace of $L_{p}(0,1)$ on which there is a linear projection $P$ of norm $\lambda$. We are interested to embed $X$ in $l_{p}^{N}$ with $N$ small so that on the image of $X$ in $l_{p}^{N}$ there will also be a projection with essentially the same norm. For $X=l_{2}^{n}, 1<p<\infty$, it is easy to see and well known that we may take (ignoring $\varepsilon$ for the moment) $N \approx n^{p^{* / 2}}$ where $p^{*}=\max \left(p, p^{\prime}\right), p^{\prime}$ being the conjugate exponent to $p$. This $N$ is large enough to ensure that $l_{2}^{n}$ embeds in both $l_{p}^{N}$ and $l_{p^{\prime}}^{N}$ (the complementation is obtained in this case automatically, any inner product space contained in $L_{r}(0,1)$ with $2 \leqslant r<\infty$ is well complemented). In this section we prove, by using the empirical method, that a similar result holds for a general $n$-dimensional subspace of $L_{p}(0,1)$ (though, in the general case, complementation is no longer automatic). The first result in this direction was proved in [Sche.2]; it involved a special family of subspaces of $L_{p}(0,1)$ (the so called $X_{p}^{\omega}$ spaces). The estimate obtained in [Sche.2] does not, however, imply Theorem 8.1 below even in this special case.

In contrast to the situation for subspaces of $l_{p}^{m}$ of which there are known many interesting examples (especially if $1 \leqslant p<2$ ) the complemented subspaces of $l_{p}^{n}$ are a rather restrictive class. In some (non-precise) sense they seem to be just 'mixtures" of $l_{p}^{m}$ spaces and inner product spaces. Therefore, while the main interest of the results of the preceding sections is in the existence of economic embeddings of concrete examples, the main application of the result in the complemented case seems to be in analyzing further the structure of complemented subspaces of $l_{p}^{n}$. In some instances it is possible to prove, using Theorem 8.1 below, that a complemented subspace of $l_{p}^{n}$ has to be isomorphic to $l_{p}^{m}$ for a suitable $m$. The second part of this section contains a specific result of this nature.

Theorem 8.1. Let $X$ be an n-dimensional subspace of $L_{p}(\Omega, \mu)$ for some probability space $(\Omega, \mu)$ and $1<p<\infty$. Assume that there is a bounded linear projection $P$ from $L_{p}(\Omega, \mu)$ onto $X$ with $\|P\|=\lambda$. Then for $0<\varepsilon<(c(p) \lambda)^{-1}$ and

$$
\begin{equation*}
N=\left[c(p, \varepsilon) n^{\left.p^{p^{*} / 2}(\log n)^{2 p^{*}-1}\right]}\right. \tag{8.1}
\end{equation*}
$$

where $p^{*}=\max \left(p, p^{\prime}\right)$, there is a subspace $\tilde{X}$ of $l_{p}^{N}$ so that $d(X, \tilde{X}) \leqslant 1+\varepsilon$ and so that there is a projection $Q$ from $l_{p}^{N}$ onto $\tilde{X}$ with

$$
\begin{equation*}
\|Q\| \leqslant \lambda(1+\varepsilon c(p) \lambda) \tag{8.2}
\end{equation*}
$$

Proof. Our strategy for getting a complemented embedding of $X$ in $l_{p}^{N}$ is the following. Let $X$ and $P$ be as in the statement of the theorem and let $Y=P^{*} X^{*} \subset L_{p^{\prime}}(\mu)$.

Assume that $N_{0}$ and the points $\left\{\omega_{j}\right\}_{j=1}^{N_{0}}$ in $\Omega$ are such that

$$
\begin{gather*}
\left|N_{0}^{-1} \sum_{j=1}^{N_{0}}\right| f\left(\omega_{j}\right)^{p^{p}}-\int_{\Omega}|f|^{p} d \mu \mid<\varepsilon, \quad f \in X, \quad\|f\|_{p}=1  \tag{8.3}\\
\left.\left|N_{0}^{-1} \sum_{j=1}^{N_{0}}\right| g\left(\omega_{j}\right)\right|^{p^{\prime}}-\int_{\Omega}|g|^{p^{\prime}} d \mu \mid<\varepsilon, \quad g \in Y, \quad\|g\|_{p^{\prime}}=1  \tag{8.4}\\
\left|N_{0}^{-1} \sum_{j=1}^{N_{0}} f\left(\omega_{j}\right) g\left(\omega_{j}\right)-\int_{\Omega} f g d \mu\right|<\varepsilon, \quad f \in X, \quad g \in Y, \quad\|f\|_{p}=\|g\|_{p^{\prime}}=1 \tag{8.5}
\end{gather*}
$$

Let $T: X \rightarrow I_{p}^{N_{0}}$ be defined by $T f=\left(f\left(\omega_{1}\right), \ldots, f\left(\omega_{N_{0}}\right)\right.$ (it is convenient in our context to consider $l_{p}^{N_{0}}$ as an $L_{p}$ space over the set $\left\{1,2, \ldots, N_{0}\right\}$ with each point having measure $N_{0}^{-1}$ ) and let $S: Y \rightarrow l_{p}^{N_{0}}$ be defined similarly. Let $i: X \rightarrow L_{p}(\mu)$ be the natural embedding. Then $P i=I_{X}=$ the identity operator of $X$, thus $i^{*} P^{*}=i_{Y}^{*} \cdot P^{*}=I_{X^{*}}$ and $P^{* *}\left(i_{\mid Y}^{*}\right)^{*}=I_{X}$ (we identify canonically $X$ with $X^{* *}$ ). The requirement (8.5) means that

$$
|\langle T f, S g\rangle-\langle f, g\rangle|<\varepsilon, \quad f \in X, \quad g \in Y, \quad\|f\|_{p}=\|g\|_{p^{\prime}}=1
$$

i.e., that $\left\|S^{*} T-\left(i_{Y}^{*}\right)^{*}\right\| \leqslant \varepsilon$. Hence

$$
\left\|P^{* *} S^{*} T-I_{X}\right\| \leqslant\|P\| \cdot \varepsilon=\lambda \varepsilon .
$$

It follows that

$$
Q=T\left(P^{* *} S^{*} T\right)^{-1} P^{* *} S^{*}
$$

is a projection of norm at most $\lambda(1+\varepsilon)^{2} /(1-\lambda \varepsilon)$ from $l_{p}^{N_{0}}$ on $T X$. Since $d(T X, X) \leqslant 1+\varepsilon$ this will prove the theorem once we verify that we can take as $N_{0}$ the number appearing in the right hand side of (8.1).

It is clear, in view of Lemma 2.5, that for the argument above it suffices to require that (8.3)-(8.5) hold just for all $x$ and $y$ belonging to an $\varepsilon$-net in the surface of the unit balls of $X$, respectively $Y$.

We shall use the procedure described above for obtaining a complemented embedding twice. In the first step we shall use the method of section 2 to embed $X$ complementably in some ${t_{p}}_{N_{\mathrm{i}}}$ where $N_{1}$ is a polynomial in $n$ which is however larger than the right hand side of (8.1). This step is needed in order to insure that for $X$ (resp. $Y$ ) we have at our disposal the sharp entropy estimates (7.9) and (7.10). Once we have these estimates we use the arguments of section 7 in order to obtain a complemented embedding into $l_{p}^{N}$ for an $N$ satisfying (8.1).

We examine the first step. In order to be specific, we assume that $1<p<2$ and then $p^{*}=p^{\prime}>2$. For performing it we have first to make a change of density procedure so as to get good $L_{\infty}$ estimates for the elements in the unit balls of $X$ and $Y=P^{*} X^{*}$. By the Lemma 7.1 there is a function $F \geqslant 0$ on $\Omega$ with $\|F\|_{L_{p}(\mu)}=1$ so that if $d \nu=F^{p} d \mu$ and $J_{F}: X \rightarrow L_{p}(v)$ is defined by $J_{F} x=F^{-1} x$ then $J_{F}$ is an isometry and $\|f\|_{\infty} \leqslant n^{1 / p}\|f\|_{p}, f \in J_{F} X$. As we have already done in previous sections, we replace $F$ by $G_{1}=\left(\left(F^{p}+1\right) / 2\right)^{1 / p}$. Then $J_{G_{1}}$ is an isometry from all of $L_{p}(\mu)$ onto $L_{p}\left(\sigma_{1}\right)$ where $d \sigma_{1}=G_{1}^{p} d \mu$ and

$$
\begin{equation*}
G_{1}(\omega) \geqslant 2^{-1 / p}, \quad \omega \in \Omega ; \quad\|f\|_{\infty} \leqslant(2 n)^{1 / p}\|f\|_{p}, \quad f \in J_{G_{1}} X \tag{8.6}
\end{equation*}
$$

Note that $P_{1}=J_{G_{1}} P J_{G_{1}}^{-1}$ is a projection of norm $\lambda$ from $L_{p}\left(\sigma_{1}\right)$ onto $J_{G_{1}} X$. Under this transformation the space $Y$ is replaced by $P_{1}^{*}\left(J_{G_{1}} X\right)^{*}=J_{G_{1}^{p-1}} Y$. By applying Lemma 7.1 to the subspace $J_{G_{1}^{p-1}} Y$ of $L_{p^{\prime}}\left(\sigma_{1}\right)$ we get in a similar way a function $G_{2}$ on $\Omega$ so that

$$
\begin{equation*}
\left\|G_{2}\right\|_{L_{p}\left(\sigma_{1}\right)}=\left\|G_{2}^{p-1}\right\|_{L_{p^{\prime}}\left(\sigma_{1}\right)}=1 \tag{8.7}
\end{equation*}
$$

and so that $J_{G_{2}}: L_{p}\left(\sigma_{1}\right) \rightarrow L_{p}\left(\sigma_{2}\right)$ is an isometry onto where $d \sigma_{2}=G_{2}^{p} d \sigma_{1}$. Moreover, we have

$$
\begin{equation*}
G_{2}(\omega) \geqslant 2^{-1 / p}, \quad \omega \in \Omega, \quad\|g\|_{\infty} \leqslant(2 n)^{1 / 2}\|g\|_{p^{\prime}} \tag{8.8}
\end{equation*}
$$

for $g \in J_{G_{2}^{p-1}}\left(J_{G_{1}^{p-1}} Y\right)=J_{\left(G_{2} G_{1}\right)^{p-1}} Y \subset L_{p^{\prime}}\left(\sigma_{2}\right)$. Also for $f=J_{G_{2} G_{1}} x \in J_{G_{2} G_{1}} X$ we get

$$
\begin{align*}
\|f\|_{\infty} & \leqslant\left\|G_{2}^{-1}\right\|_{\infty}\left\|J_{G_{1}} x\right\|_{\infty} \leqslant\left\|G_{2}^{-1}\right\|_{\infty}(2 n)^{1 / p}\left\|J_{G_{1}} x\right\|_{L_{p}\left(\sigma_{1}\right)}  \tag{8.9}\\
& =\left\|G_{2}^{-1}\right\|_{\infty}(2 n)^{1 / p}\|f\|_{L_{p}\left(\sigma_{2}\right)} \leqslant(4 n)^{1 / p}\|f\|_{L_{p}\left(\sigma_{2}\right)} .
\end{align*}
$$

In conclusion, we see that by replacing $X$ and $Y$ respectively by $J_{G_{2} G_{1}} X$ and $J_{\left(G_{2} G_{1}\right)^{p-1}} Y$ there is no loss of generality to assume that for $f \in X$ with $\|f\|_{p}=1$ and $g \in Y$ with $\|g\|_{p^{\prime}}=1$

$$
\begin{equation*}
\left\|f^{p}\right\|_{\infty} \leqslant 4 n, \quad\left\|g^{p^{\prime}}\right\|_{\infty} \leqslant(4 n)^{p^{\prime} / 2}, \quad\|f g\|_{\infty} \leqslant(4 n)^{1 / 2+1 / p} \tag{8.10}
\end{equation*}
$$

By (8.10) and the method of section 2 we get that (8.3)-(8.5) hold for suitable $\left\{\omega_{j}\right\}_{j=1}^{N_{0}}$ if

$$
N_{0} \geqslant c_{1} n \varepsilon^{-2} \log \varepsilon^{-1} \cdot(4 n)^{p^{\prime / 2}}
$$

We now apply Proposition 7.2 to obtain, via another change of density, good estimates on the entropy numbers. Since we work with densities which are bounded from below an argument similar to the one we did above (and using the definition of the entropy numbers) shows that we can make a change of density which is good for both $X$ and $Y$. In other words there is no loss of generality to assume that for suitable constants $c_{2}$ and $c_{3}(p)$ we have

$$
\begin{align*}
& \|f\|_{\infty} \leqslant c_{2} n^{1 / p}\|f\|_{p}, \quad f \in X ; \quad\|g\|_{\infty} \leqslant c_{2} n^{1 / 2}\|g\|_{p^{\prime}}, \quad g \in Y  \tag{8.11}\\
& \log E\left(B_{X}, B_{\infty}, t\right) \leqslant c_{3}(p) \log n \varepsilon^{-1} \cdot n t^{-p}, \quad 1 \leqslant t \leqslant c_{2} n^{1 / p}  \tag{8.12}\\
& \log E\left(B_{Y}, B_{\infty}, t\right) \leqslant c_{3}(p) \log n \varepsilon^{-1} \cdot n t^{-2}, \quad 1 \leqslant t \leqslant c_{2} n^{1 / 2} \tag{8.12}
\end{align*}
$$

From the proof of Theorem 7.3 it follows now that if we take as $N_{0}$ the right hand side of (8.1) then (8.3) and (8.4) hold for suitable $\left\{\omega_{j}\right\}_{j=1}^{N_{0}}$ in $\Omega$ and all $x$, resp. $y$, in a suitable $\varepsilon$-net $\mathscr{F}_{X}$, resp. $\mathscr{F}_{Y}$, in the boundary of the unit balls of $X$, resp. $Y$. By an argument which is similar to the proof of Theorem 7.3 we shall show next that also (8.5) holds for this $N_{0}$ and this will conclude the proof.

As in the proof of Theorem 7.3, we can replace the functions $x \in \mathscr{F}_{X}$ and $y \in \mathscr{F}_{Y}$ by the functions of the form

$$
\begin{aligned}
& \tilde{x}=x \chi_{D_{0, x}}+\sum_{k=1}^{l_{1}}(1+\varepsilon)^{k}\left(\chi_{D_{k, x}^{\prime}}-\chi_{D_{k, x}^{\prime \prime}}\right)=\sum_{k=0}^{l_{1}} \varphi_{k, x} \\
& \tilde{y}=y \chi_{D_{0, y}}+\sum_{h=1}^{l_{2}}(1+\varepsilon)^{h}\left(\chi_{D_{h, y}^{\prime}}-\chi_{D_{h, y}^{\prime \prime}}\right)=\sum_{h=0}^{l_{2}} \psi_{h, y}
\end{aligned}
$$

where

$$
l_{1}=\left[\log \left(c_{2} n^{1 / p}\right) / \log (1+\varepsilon)\right], \quad l_{2}=\left[\log \left(c_{2} n^{1 / 2}\right) / \log (1+\varepsilon)\right]
$$

and for each $x \in \mathscr{F}_{X}$, the sets $D_{0, x}, D_{k, x}^{\prime}, D_{k, x}^{\prime \prime}, 1 \leqslant k \leqslant l_{1}$ form a decomposition of $\Omega$ into disjoint sets, a similar statement holding for $y \in \mathscr{F}_{Y}$. In the present context we have to take into consideration also the signs of $x$ and $y$ and this is the reason for the appearance of $D_{k, x}^{\prime}$ and $D_{k, x}^{\prime \prime}$ rather than just $D_{k, x}$ as in the proof of 7.3.

As in the proof of 7.3 (see (7.26)) we get from the construction and (8.12) and (8.13) that if we put $\mathscr{B}_{k}=\left\{\varphi_{k, x}\right\}_{x \in \mathscr{F}_{X}}$ and $\mathscr{C}_{h}=\left\{\psi_{h, y}\right\}_{y \in \mathscr{F}_{r}}$ then

$$
\begin{equation*}
\log \overline{\mathscr{B}}_{k} \leqslant c_{4}(p, \varepsilon) n \log n /(1+\varepsilon)^{k p}, \quad \log \overline{\bar{C}}_{h} \leqslant c_{4}(p, \varepsilon) n \log n /(1+\varepsilon)^{2 h} \tag{8.14}
\end{equation*}
$$

We consider the maps $T$ and $S$ defined after (8.5) with $N_{0}$ being the $N$ given in (8.1). (Obviously $T$ and $S$ are defined not just on $X$ and $Y$ but on all the functions on $\Omega$.) As mentioned above the argument of the proof of Theorem 7.3 and (8.1) ensure that we can assume that

$$
\sum_{k=0}^{l_{1}}\left|\left\|\varphi_{k, x}\right\|_{L_{p}(\Omega)}^{P_{\Omega}}-\left\|T \varphi_{k, x}\right\|_{l_{p}^{N}}^{p_{p}}\right| \leqslant \varepsilon, \quad x \in \mathscr{F}_{X}
$$

and that a similar statements holds for $y \in \mathscr{F}_{Y}$. Our goal is to show that we can ensure also that

$$
\begin{equation*}
|\langle\boldsymbol{T} \tilde{x}, S \tilde{y}\rangle-\langle\tilde{x}, \tilde{y}\rangle| \leqslant \sum_{k=0}^{l_{1}} \sum_{h=0}^{l_{2}}\left|\left\langle\boldsymbol{T} \varphi_{k, x}, S \psi_{h, y}\right\rangle-\left\langle\varphi_{k, x}, \psi_{h, y}\right\rangle\right|<\varepsilon \tag{8.15}
\end{equation*}
$$

for all $x \in \mathscr{F}_{X}$ and $y \in \mathscr{F}_{Y}$.
Note that for $k \geqslant 1,\left|\varphi_{k, x}\right|$ is $(1+\varepsilon)^{k}$ times the characteristic function of a set. Because of the special form of $T,\left|T \varphi_{k, x}\right|$ is also $(1+\varepsilon)^{k}$ times the characteristic function of a set with essentially the same measure. The same holds for $\left|\psi_{h, y}\right|$ and $\left|S \psi_{h, y}\right|$. Our strategy for verifying (8.15) is to show that the terms in the sum of $(8.15)$ for which $k p$ is far from $h p^{\prime}$ are small because $\left|\left\langle\varphi_{k, x}, \psi_{h, y}\right\rangle\right|$ are small for these indices, while the rest of the terms can be handled by the estimates provided by the empirical distribution theory (i.e., Lemma 2.1). In view of (2.4) and (8.14) the $h$ and $k$ for which we can apply the empirical distribution estimate are those for which

$$
\begin{equation*}
N /(1+\varepsilon)^{h+k}>c_{5}(p, \varepsilon) n \log n \max \left((1+\varepsilon)^{-k p},(1+\varepsilon)^{-2 h}\right) . \tag{8.16}
\end{equation*}
$$

Let $\alpha_{n}$ be such that

$$
(1+\varepsilon)^{\alpha_{n}}=(\log n)^{2 p /(p-1)},
$$

and consider $h$ and $k$ for which

$$
\begin{equation*}
2 h<k p<h p^{\prime}+\alpha_{n} . \tag{8.17}
\end{equation*}
$$

In this case (8.16) becomes

$$
N>c_{5}(p, \varepsilon) n \log n(1+\varepsilon)^{k-h}
$$

but

$$
(1+\varepsilon)^{k-h}<(1+\varepsilon \varepsilon)^{h\left(p^{\prime} / p-1\right)}(1+\varepsilon)^{\alpha_{n}^{\prime} p} \leqslant\left(c_{2} n^{1 / 2}\right)^{p^{\prime / p-1}} \log n^{2(p-1)}
$$

and thus since $1+\frac{1}{2}\left(p^{\prime} / p-1\right)=p^{\prime} / 2,(8.1)$ ensures that (8.16) holds for indices satisfying (8.17). Another case for which (8.16) is easily seen to hold is

$$
\begin{equation*}
k p<2 h, \quad h \leqslant \log \log n \tag{8.18}
\end{equation*}
$$

For the $h, k$ for which (8.17) and (8.18) fail we will show that the terms $\left|\left\langle\varphi_{k, x}, \psi_{h, y}\right\rangle\right|$ contribute only a negligible amount to (8.15) (and similarly $\left|\left\langle T \varphi_{k, x}, S \psi_{h, y}\right\rangle\right|$ ). Recall that $\int_{\Omega}\left|\varphi_{k, x}\right|^{p} d \nu \leqslant 2$ and $\int_{\Omega}\left|\psi_{h, y}\right|^{p^{\prime}} d \nu \leqslant 2$, hence for all $x, y, k$ and $h$

$$
\begin{equation*}
\left|\left\langle\varphi_{k, x}, \psi_{h, y}\right\rangle\right| \leqslant 2(1+\varepsilon)^{h+k} \min \left((1+\varepsilon)^{-k p},(1+\varepsilon)^{-h p^{\prime}}\right) \tag{8.19}
\end{equation*}
$$

Hence

$$
\begin{align*}
\sum_{k p>h p^{\prime}+\alpha_{n}}\left|\left\langle\varphi_{k, x}, \psi_{h, y}\right\rangle\right| & \leqslant c_{6}(p, \varepsilon) \sum_{h=0}^{l_{1}}(1+\varepsilon)^{h+\left(h p^{\prime}+\alpha_{n}\right)(1-p) / p} \\
& =c_{6}(p, \varepsilon) \sum_{h=0}^{l_{1}}(1+\varepsilon)^{\alpha_{n}(1-p) / p}  \tag{8.20}\\
& =c_{6}(p, \varepsilon) l_{1} /(\log n)^{2} \leqslant c_{7}(p, \varepsilon) / \log n .
\end{align*}
$$

We also have

$$
\begin{align*}
\sum_{\substack{k p<2 h \\
h>\log \log n}}\left|\left\langle\varphi_{k, x}, \psi_{h, y}\right\rangle\right| & \leqslant 2 \varepsilon^{-1} \sum_{h=\log \log n}^{l_{1}}(1+\varepsilon)^{h\left(2-p^{\prime}\right) / p}  \tag{8.21}\\
& \leqslant 2 \varepsilon^{-1}(\log n)^{-c_{8}(p, \varepsilon)}
\end{align*}
$$

for some $c_{8}(p, \varepsilon)>0$.
(8.20) and (8.21) complete the verification of (8.15) and the proof is finally concluded.

We apply Theorem 8.1. to prove a result in the spirit of [B.Tz.1].

Proposition 8.2. Assume that $l_{p}^{n}=X \oplus Y$ with $m=\operatorname{dim} Y<n^{2 / p^{*}-\varepsilon}$ for some $1<p<\infty$ and $\varepsilon>0$ (as always $p^{*}=\max \left(p, p^{\prime}\right)$ ). Then

$$
\begin{equation*}
d\left(X, l_{p}^{n-m}\right) \leqslant c(\varepsilon, \lambda, M, p) \tag{8.22}
\end{equation*}
$$

where $\lambda$ is norm of the projection of $l_{p}^{n}$ on $Y$ which maps $X$ to 0 and $M$ is the basis constant of $Y$ (i.e., the inf of the basis constants of all possible bases of $Y$ ).

Before we prove Proposition 8.2. we state an immediate corollary of it for spaces arising in harmonic analysis. Recall that for a subset $\Lambda$ of the integers $L_{p, \mathrm{~A}}$ is the subspace of $L_{p}[0,2 \pi]$ of all the functions whose Fourier coefficients vanish outside $\Lambda$.

Proposition 8.3. Let $1<p<\infty, \varepsilon>0$ and $\Lambda \subset\{1, \ldots, n\}$ satisfy $\overline{\bar{\Lambda}}<n^{2 / p^{*-\varepsilon}}$. Then

$$
\begin{equation*}
d\left(L_{p,\{1, \ldots, n\} \backslash \Lambda}, l_{p}^{n-\tilde{\Lambda}}\right) \leqslant c(\varepsilon, \lambda, p), \tag{8.23}
\end{equation*}
$$

where $\lambda$ is the norm of the orthogonal projection from $L_{p}(0,2 \pi)$ onto $L_{p, \Lambda}$.

Proof. By a classical result $d\left(L_{p,\{1, \ldots, n\}}, l_{p}^{m}\right) \leqslant c_{1}(p)$ (see [Zy], Chapter X, Theorem (7.5)). Another classical fact, a theorem of M. Riesz, asserts that the characters in their natural order form a Schauder basis of $L_{p}(0,2 \pi)$. Hence the basis constant of $L_{p, \Lambda}$ is bounded by a number depending only on $p$. These two classical facts show that (8.23) is a consequence of (8.22).

We turn to the proof of Proposition 8.2 which follows some arguments of [B.Tz.1]. The proof consists of two steps.

Step 1. For fixed $Y$ with $\operatorname{dim} Y=m<n^{2 / p^{*}-\varepsilon}$ the relation $l_{p}^{r}=Y \oplus X$ has a "unique" solution. In other words, if $\eta_{i}=d\left(l_{p}^{n}, Y \oplus X_{i}\right), i=1,2$, then $d\left(X_{1}, X_{2}\right) \leqslant C\left(\varepsilon, p, \lambda_{1}, \lambda_{2}, \eta_{1}, \eta_{2}\right)$ where $\lambda_{1}$, resp. $\lambda_{2}$, is the norm of the projection onto $Y$ taking $X_{1}$, resp. $X_{2}$, to 0 . In this step the basis constant of $Y$ plays no role.

Step 2. With $Y$ as in the statement of Proposition 8.2.

$$
d\left(l_{p}^{n}, l_{p}^{n-m} \oplus Y\right) \leqslant C_{2}(\varepsilon, \lambda, M, p)
$$

where the direct sum in $l_{p}^{n-m} \oplus Y$ is taken, e.g., in the $l_{p}$ sense.
Proof of Step 1. In the sequel, all the isomorphism and complementation constants are assumed to be controlled by the parameters of the given data. To prove that the solution $X$ of $l_{p}^{n} \approx X \oplus Y$ is unique, it suffices by [B.Tz.1] to prove that $X$ contains a subspace $Z$ complemented in $l_{p}^{n}$, so that $\operatorname{dim} Z=n^{\prime}>c n$, that both $Z$ and its complement are isomorphic to $l_{p}$ spaces of the appropriate dimension and that moreover $Y$ embeds complementably in $l_{p}^{r^{\prime}}$. The latter requirement follows from Theorem 8.1 since $n^{\prime}>m^{p^{*} / 2}(\log m)^{2 p^{*}-1}$. The existence of $Z$ follows from [J.Sche.2] or [B.Tz.2]. Indeed, let
$P$ be the projection on $Y$ whose kernel is $X$. By [B.Tz.2] there is a subset $A \subset\{1,2, \ldots, n\}$ with $\overline{\bar{A}} \geqslant c n$ so that

$$
\begin{equation*}
\left\|R_{A} P R_{A}\right\|_{p \rightarrow p}<\frac{1}{10} \tag{8.24}
\end{equation*}
$$

where $R^{A}$ is the natural projection from $l_{p}^{n}$ onto span $\left\{e_{i}\right\}_{i \in A}$. Let

$$
Z=(I-P) \operatorname{span}\left\{e_{i}\right\}_{i \in A} \subset X
$$

Clearly by (8.24), $l_{p}^{n}=Z \oplus \operatorname{span}\left\{e_{i}\right\}_{i ₫ A}$.
Proof of Step 2. We shall use the finite-dimensional version of the decomposition method introduced in [B.D.G.J.N.]. Assume that for some $m_{0}$ and every subspace $Y$ of dimension $m_{0}$ of $l_{p}$ so that $Y$ is complemented in $l_{p}$ and has a good basis, we know that for some $n_{0}$

$$
\begin{equation*}
Y \oplus l_{p}^{n_{0}} \approx l_{p}^{n_{0}+m_{0}} \tag{8.25}
\end{equation*}
$$

with a good control of the constants. We shall verify that then for every $m_{1}$ a similar statement is valid if $m_{0}$ and $n_{0}$ are replaced by $m_{0} m_{1}$ and $m_{1}\left(\left(m_{0} m_{1}\right)^{p^{* / 2}+\varepsilon}+n_{0}\right)$ respectively, again with a good control of the constants. This argument will show that if $\varrho$ is defined by $n_{0}=m_{0}^{\rho}$ and if $\varrho>p^{* / 2}$ then (by an appropriate choice of $m_{1}$ ) we can replace $\varrho$ by any $\tilde{\varrho}$ satisfying

$$
\varrho>\frac{p^{*}}{2}\left(1-\frac{1}{\varrho}+\frac{2}{p^{*}}\right) .
$$

In other words we will get that (8.25) is valid for every $m_{0}$ and $n_{0}=m_{0}^{\varrho}$ provided $\varrho>p^{* / 2}$ and this is the assertion in Step 2.

Let now $m=m_{0} m_{1}$ and $n=m^{p^{*} / 2+\varepsilon}$. Let $Y$ be a complemented subspace of $l_{p}$ with a Schauder basis $\left\{u_{i}\right\}_{i=1}^{m}$. Write $\{1, \ldots, m\}$ as a union of $m_{1}$ consecutive intervals $A_{j}$ each of size $m_{0}$ and put

$$
W_{j}=\operatorname{span}\left\{u_{i}\right\}_{i \in A_{j^{\prime}}} \quad j=1, \ldots, m_{1}
$$

By the assumption (8.25) we have for every $j$

$$
\begin{equation*}
W_{j} \oplus l_{p}^{n_{0}} \approx l_{p}^{n_{0}+m_{0}} \tag{8.26}
\end{equation*}
$$

Let

$$
U_{j}=\operatorname{span}\left\{u_{i}\right\}_{i \in U_{k \geqslant j} A_{k}}, j=1, \ldots, m_{1}
$$

Note that $U_{1}=Y$ and that $U_{m_{1}}=W_{m_{1}}$. Since $\operatorname{dim} U_{j} \leqslant m$ for every $j$ it follows by our choice of $n$ and Theorem 8.1 that for every $j$ and suitable $Z_{j}$

$$
\begin{equation*}
l_{p}^{n} \approx U_{j} \oplus Z_{j} . \tag{8.27}
\end{equation*}
$$

By (8.26) and (8.27) we get

$$
\begin{equation*}
l_{p}^{n+n_{0}} \approx U_{j} \oplus Z_{j} \oplus l_{p}^{n_{0}} \approx U_{j+1} \oplus W_{j} \oplus Z_{j} \oplus l_{p}^{n_{0}} \approx U_{j+1} \oplus Z_{j} \oplus l_{p}^{n_{0}+m_{0}} . \tag{8.28}
\end{equation*}
$$

Write

$$
l_{p}^{\left.m_{1}-1\right) n+m_{1}\left(m_{0}+n_{0}\right)}=I_{p}^{n+m_{0}+n_{0}} \oplus \ldots \oplus I_{p}^{n+m_{0}+n_{0}} \oplus l_{p}^{n_{0}+m_{0}}
$$

with $m_{1}$ summands (the direct sums are in the $l_{p}$ sense). By (8.26), (8.27) and (8.28) we get

$$
\begin{aligned}
l^{\left(m_{1}-1\right) n+m_{1}\left(m_{0}+n_{0}\right)} & \approx\left(U_{1} \oplus Z_{1} \oplus l_{p}^{n_{0}+m_{0}}\right) \oplus \ldots \oplus\left(U_{m_{1}-1} \oplus Z_{m_{1}-1} \oplus l_{p}^{n_{0}+m_{0}}\right) \oplus\left(W_{m_{1}} \oplus l_{p}^{n_{0}}\right) \\
& \approx U_{1} \oplus\left(U_{2} \oplus Z_{1} \oplus l_{p}^{n_{0}+m_{0}}\right) \oplus \ldots \oplus\left(U_{m_{1}} \oplus Z_{m_{1}-1} \oplus l_{p}^{n_{0}+m_{0}}\right) \oplus l_{p}^{n_{0}} \\
& \approx U_{1} \oplus l_{p}^{n+n_{0}} \oplus \ldots \oplus l_{p}^{n+m_{0}} \oplus l_{p}^{n+m_{0}} \oplus l_{p}^{n_{0}} \\
& =Y \oplus l_{p}^{\left(m_{1}-1\right) n+m_{1} n_{0}}
\end{aligned}
$$

and this concludes the proof.

## 9. Entropy estimates, sharper results

The number of balls of radius $k$ in $l_{\infty}^{n}$ needed to cover the ball of radius $n$ in $l_{1}^{n}$ (=the unit ball of $L_{1}(\mu)$ where $\mu$ is the natural probability measure on $\left.(1,2, \ldots, n\}\right)$ is clearly equivalent to the number of solutions to the equation

$$
\left|x_{1}\right|+\left|x_{2}\right|+\ldots+\left|x_{n}\right| \leqslant n
$$

where each $x_{i}$ is an integer (positive, negative or 0 ) divisible by $k$. This number can be computed explicitely and as observed in [Schu.] one gets from this computation that

$$
\begin{equation*}
\log E\left(n B_{l_{1}^{n}}, B_{l_{m}^{n}}, t\right) \approx n(\log t) / t, \quad 2 \leqslant t \leqslant n . \tag{9.1}
\end{equation*}
$$

In this section we shall prove a generalization of (9.1) to the case where $l_{1}^{n}$ is replaced by an $n$-dimensional subspace of $l_{1}^{N}$ in a "good position". The result we obtain
will be sharp and strengthen the results obtained in section 4. Like (9.1) the result has a natural interpretation as an estimate for the cardinality of a suitable set of lattice points. The counting of the number of lattice points in the ball of an arbitrary $n$-dimensional subspace $X$ of $l_{1}^{N}$ is naturally more complicated than in the concrete case of $l_{1}^{n}$. The results in this section are obtained by refining the Banach space techniques of section 4. Hence one may view the results of this section as an application of Banach space theory to a problem of counting lattice points.

As in section 4 the estimate of the $L_{\infty}$ entropy of the $L_{1}$ ball will be done via $L_{2}$ estimates. Thus our problem naturally divides into two parts-estimating the $L_{2}$ entropy of the $L_{1}$ ball and the $L_{\infty}$ entropy of the $L_{2}$ ball. We start with the second part.

Proposition 9.1. Let $X$ be an n-dimensional subspace of $L_{1}(\mu)$ where $\mu$ is a probability measure on $\{1, \ldots, N\}$ satisfying $\mu\{i\} \geqslant 1 / 2 N$ for every $i$. Denote by $X_{2}$ resp. $X_{\infty}$ the space $X$ endowed with the norm induced by $L_{2}(\mu)$ resp. $L_{\infty}(\mu)$. Let $\left\{\varphi_{j}\right\}_{j=1}^{n}$ be an orthonormal basis of $X_{2}$ and define $\varrho$ by

$$
\begin{equation*}
\left\|\left(\sum_{j=1}^{n} \varphi_{j}^{2}\right)^{1 / 2}\right\|_{\infty}=\varrho n^{1 / 2} \tag{9.2}
\end{equation*}
$$

Then

$$
\begin{equation*}
\log E\left(B_{X_{2}}, B_{X_{\infty}}, t\right) \leqslant c \varrho^{2} n t^{-2} \log (t N / n), \quad 2 \leqslant t \leqslant \varrho n^{1 / 2} \tag{9.3}
\end{equation*}
$$

where $c$ is an absolute constant.

We prove (9.3) by using another intermediate space, namely the Orlicz space $L_{\psi_{2}}(\mu)$ where (as in section 6) $\psi_{2}(s)=\exp \left(s^{2}\right)-1$.

Lemma 9.2. With $X$ as in Proposition 9.1 we have for a suitable absolute constant $c_{1}$

$$
\begin{equation*}
\log E\left(B_{X_{2}}, B_{X_{\psi_{2}}}, t\right) \leqslant c_{2} \varrho^{2} n t^{-2}, \quad 1 \leqslant t \leqslant \varrho n^{1 / 2} \tag{9.4}
\end{equation*}
$$

Proof. By Proposition 4.2 we have

$$
\begin{equation*}
\log E\left(B_{X_{2}}, B_{X_{\psi_{2}}}, t\right) \leqslant c_{2} n t^{-2} M^{2} \tag{9.5}
\end{equation*}
$$

where

$$
M=n^{-1 / 2} \int_{\Omega}\left\|\sum_{j=1}^{n} g_{j}(\omega) \varphi_{j}\right\|_{L_{v_{2}}(\omega)} d \sigma(\omega)
$$

the $\left\{g_{j}\right\}_{j=1}^{n}$ being normalized independent Gaussians on the probability space $(\Omega, \sigma)$.
Observe that $\|f\|_{L_{\nu_{2}}(\omega)} \leqslant c_{3} \int \exp \left(f^{2}\right) d \mu$ for every function $f$ and some universal $c_{3}$. Hence for every $c_{4}>0$ there is a $c_{5}<\infty$ so that

$$
\begin{aligned}
\varrho^{-1} n^{-1 / 2} \int_{\Omega}\left\|\sum_{j=1}^{n} g_{j}(\omega) \varphi_{j}\right\|_{L_{\psi_{2}}(\mu)} d \sigma(\omega) & \leqslant c_{5} \int_{\Omega}\left(\int \exp \left(c_{4}\left|\sum_{j=1}^{n} g_{j}(\omega) \varphi_{j}\right|^{2} / n \varrho^{2}\right) d \mu\right) d \sigma(\omega) \\
& =c_{5} \int\left(\int_{\Omega} \exp \left(c_{4}\left|\sum_{j=1}^{n} g_{j}(\omega) \varphi_{j}\right|^{2} / n \varrho^{2}\right) d \sigma(\omega)\right) d \mu .
\end{aligned}
$$

By the Laudau-Schepp-Fernique theorem (see [Ma.P.]) we have that for suitable $c_{4}$ and $c_{6}$

$$
\int_{\Omega} \exp \left(c_{4}\left|\sum_{j=1}^{n} g_{j}(\omega) \alpha_{j}\right|^{2}\right) d \sigma(\omega) \leqslant c_{6}
$$

whenever $\sum_{j=1}^{n} a_{j}^{2} \leqslant 1$. Consequently, we deduce that $M \leqslant \varrho c_{5} c_{6}$ and by substituting this estimate in (9.5) we get (9.4).

Lemma 9.3 With $\mu$ as in the statement of Proposition 9.1 we have

$$
\begin{equation*}
\log E\left(B_{L_{\psi_{2}}(\mu)}, B_{L_{x_{x}}(\mu)}, t\right) \leqslant c_{7} \exp \left(-c_{8} t^{2}\right) N, \quad 1 \leqslant t \leqslant(\log N)^{1 / 2} . \tag{9.6}
\end{equation*}
$$

Proof. Let $f \in B_{L_{\psi_{2}}(\omega)}$ and put $A=\{i ;|f(i)| \leqslant 2 t\}, f_{0}=f \chi_{A}, f_{1}=f-f_{0}$. A direct computation shows that $\left\|f_{1}\right\|_{L_{2}(\mu)} \leqslant c_{9} \exp \left(-t^{2}\right)$. Hence

$$
\begin{equation*}
B_{L_{\psi_{2}}(\mu)} \subset 2 t B_{L_{L_{x}}(\mu)}+c_{9} \exp \left(-t^{2}\right) B_{L_{2}(\mu)} . \tag{9.7}
\end{equation*}
$$

Since $\mu\{i\}>1 / 2 N$ for every $i$ we get from (9.1) (cf. [Schu.]) that

$$
\log E\left(B_{L_{2}(\mu)}, B_{L_{\alpha}(\mu)}, s\right) \leqslant c_{10} N s^{-2} \log s \leqslant c_{11} N s^{-1}
$$

and hence, by (9.7),

$$
\log E\left(B_{L_{\psi_{2}}(\mu)}, B_{L_{\infty}(\mu)}, c_{12} t\right) \leqslant c_{13} N \exp \left(-t^{2}\right)
$$

and this implies (9.6).

Proof of Proposition 9.1. For $1<s<t$ we have by Lemmas 9.2 and 9.3

$$
\begin{aligned}
\log E\left(B_{X_{2}}, B_{X_{x}}, t\right) & \leqslant \log E\left(B_{X_{2}}, B_{X_{\psi_{2}}}, t s^{-1}\right)+\log E\left(B_{L_{\psi_{2}}(\mu)}, B_{L_{\alpha}(\mu)}, s\right) \\
& \leqslant c_{1} \varrho^{2} n s^{2} t^{-2}+c_{7} N \exp \left(-c_{8} s^{2}\right) .
\end{aligned}
$$

By taking $s \approx(\log N t / n)^{1 / 2}$, (9.3) follows.

Corollary 9.4. Let $X$ be as in Proposition 9.1. and let $Y$ be an m-dimensional subspace of $X$. Then there is an $f \in Y$ satisfying

$$
\begin{equation*}
\|f\|_{L_{2}(\mu)}=1, \quad\|f\|_{L_{\alpha}(\mu)} \leqslant c_{14} \varrho\left(n m^{-1} \log \left(N \varrho m^{-1}\right)\right)^{1 / 2} \tag{9.8}
\end{equation*}
$$

Proof. Assume that for some $t>1, t B_{Y_{x}} \subset \frac{1}{2} B_{Y_{2}}$. Then by (9.3)

$$
\begin{aligned}
2 m & \leqslant \log E\left(B_{Y_{2}}, B_{Y_{x}}, t\right) \leqslant \log E\left(B_{X_{2}}, B_{X_{x}}, t\right) \\
& \leqslant c p^{2} n t^{-2} \log (t N / n),
\end{aligned}
$$

and hence $2 t<c_{14} \varrho\left(n m^{-1} \log (N \varrho / m)\right)^{1 / 2}$.

Corollary 9.5. Let $X$ be an $n$-dimensional subspace of $l_{\infty}^{N}$. Then

$$
\begin{equation*}
d\left(X, l_{2}^{n}\right) \geqslant c_{15} n^{1 / 2}(1+\log N / n)^{-1 / 2} \tag{9.9}
\end{equation*}
$$

Proof. Let $\mu$ be any probability measure on $\{1, \ldots, N\}$ satisfying $\mu\{i\}>1 / 2 N$ for every $i$ and let $\varrho$ be given by (9.2). By Corollary 9.4 there is an orthogonal system $f_{1}, \ldots, f_{m}(m=n / 2)$ in $X_{2}$ satisfying $\left\|f_{j}\right\|_{L_{x}(\mu)} \leqslant c_{16} \varrho(\log (N \varrho / n))^{1 / 2}$ for every $j$. By the orthonormality of the $f_{j}$

$$
\int\left\|\sum_{j=1}^{m} f_{j}(i) f_{j}\right\|_{X} d \mu(i) \leqslant\left(\int\left\|\sum_{j=1}^{m} f_{j}(i) f_{j}\right\|_{X}^{2} d \mu(i)\right)^{1 / 2} \leqslant m^{1 / 2} d\left(X, l_{2}^{n}\right) c_{16} \varrho(\log (N \varrho / n))^{1 / 2}
$$

On the other hand since the norm in $X$ is the $L_{\infty}$ norm

$$
\int\left\|\sum_{j=1}^{m} f_{j}(i) f_{j}\right\|_{X} d \mu(i) \geqslant \int\left(\sum_{j=1}^{m} f_{j}(i)^{2}\right) d \mu(i)=m
$$

Hence

$$
d\left(X, l_{2}^{m}\right) \geqslant m^{1 / 2}\left(c_{16} \varrho(\log (N \varrho / n))^{1 / 2}\right)^{-1}
$$

We are free to choose the measure $\mu$. We show that $\mu$ can be chosen so that $\varrho$ will be bounded by an absolute constant and this will prove (9.9). The 2 -absolutely summing norm of the identity of $X$ is $n^{1 / 2}$ (cf. e.g., [P4], Theorem 1.11) and hence by the Pietch factorization theorem there is a probability measure $v$ on $\{1, \ldots, N\}$ so that $\|x\|_{\infty} \leqslant n^{1 / 2}\|x\|_{L_{2}(v)}$ for every $x \in X$. Define $\mu$ by

$$
\mu\{i\}=\frac{1}{2 N}+\frac{1}{2} v\{i\}, \quad 1 \leqslant i \leqslant N
$$

For every orthonormal system $\left\{\varphi_{j}\right\}_{j=1}^{n}$ of $X$ (in the norm induced by $L_{2}(\mu)$ and any choice of scalars $\left\{\alpha_{j}\right\}_{j=1}^{n}$ we have

$$
\left\|\sum_{j=1}^{n} \alpha_{j} \varphi_{j}\right\|_{\infty} \leqslant 2 n^{1 / 2}\left\|\sum_{j=1}^{n} \alpha_{j} \varphi_{j}\right\|_{L_{2}(\mu)}=2 n^{1 / 2}\left(\sum_{j=1}^{N} a_{j}^{2}\right)^{1 / 2}
$$

Consequently $\left\|\left(\sum_{j=1}^{n} \varphi_{j}^{2}\right)^{1 / 2}\right\|_{\infty} \leqslant 2 n^{1 / 2}$.

Remarks. (1) S. Kislyakov informed us that E. Gluskin recently obtained (9.9).
(2) A simple example shows how precise the estimate (9.9) is (and for that matter also the preceding estimates). Let $X(k, m)=\left(\Sigma \oplus l_{2}^{k}\right)_{\infty}$ with $m$ summands in the direct sum. Then $n=\operatorname{dim} X(k, m)=k m, X(k, m)$ is 2 isomorphic to a subspace of $l_{\infty}^{N}$ with $N=m 2^{c k}$ and $d\left(X(k, m), l_{2}^{n}\right) \leqslant m^{1 / 2}$. It follows that for all choices of $k$ and $m$, (9.9) gives an equivalence.
(3) In [F.L.M.] it is shown that $d\left(X, l_{2}^{n}\right) \geqslant c n^{1 / 2}(\log N)^{-1 / 2}$ and in [F.J.] it is proved that $d\left(X, l_{2}^{\prime \prime}\right) \geqslant \frac{1}{2} n^{1 / 2}(n / N)^{1 / 2}$. Thus only if $N$ is close to $n$ but not too close (e.g., $N=n \log n$ ) is (9.9) an improvement on previously published results.
(4) B. Carl and A. Pajor [C.Pa.] obtained independently an estimate which is essentially the dual of Proposition 9.1 (and thus by [T-J.] also essentially equivalent to Proposition 9.1).

We pass to the estimate of the $L_{2}$ entropy of the unit ball of a subspace $X$ of $L_{1}$. Let $X$ be an $n$-dimensional subspace $l_{1}^{N}$. As we noted in Lemma 4.5 it follows from the lemma of Lewis that there is a probability measure $\mu$ on $\{1, \ldots, N\}$ with $\mu\{i\} \geqslant 1 / 2 N$ for every $i$, a subspace $\bar{X}$ of $L_{1}(\mu)$ which is isometric to $X$ and an orthonormal system $\left\{\varphi_{j}\right\}_{j=1}^{n}$ in $\bar{X}$ so that

$$
\begin{equation*}
\sum_{j=1}^{n} \varphi_{j}^{2}(i) \leqslant 4 n, \quad 1 \leqslant i \leqslant N \tag{9.10}
\end{equation*}
$$

The entropy estimate we are going to prove will be for the unit ball of $\bar{X}$.
For the proof below it is convenient (though not really necessary) to note that if $B$ is the unit ball of an inner product space and $D$ is any compact convex set then $E(D, B, t)=\tilde{E}(D, B, t)$. Indeed, if a compact convex set $K$ is in the unit ball it is also contained in a ball with radius 1 centered at $K$ (take as center the nearest point to 0 in $K$ ).

Proposition 9.6. Let $\hat{X}$ be an n-dimensional subspace of $L_{1}(\mu)$ as above. Then

$$
\begin{equation*}
\log E\left(B_{\bar{X}_{1}}, B_{\tilde{X}_{2}}, t\right) \leqslant c \min \left(K(\tilde{X})^{2}, \log t\right) n t^{-2}, \quad 2 \leqslant t \leqslant n^{1 / 2} \tag{9.11}
\end{equation*}
$$

where $K(\bar{X})$ is $K$-convexity constant of $\bar{X}$.

Proof. We have

$$
\begin{equation*}
\log E\left(B_{\tilde{X}_{1}}, B_{\tilde{X}_{2}}, t\right) \leqslant \log E\left(B_{\tilde{X}_{1}}, B_{\tilde{X}_{1}} \cap t B_{\tilde{X}_{2}}, 2\right)+\log E\left(B_{\tilde{X}_{1}} \cap t B_{\tilde{X}_{2}}, B_{\tilde{X}_{2}}, \frac{t}{2}\right) \tag{9.12}
\end{equation*}
$$

Note that

$$
\begin{equation*}
E\left(B_{\bar{X}_{1}}, B_{\bar{X}_{1}} \cap t B_{\tilde{X}_{2}}, 2\right) \leqslant E\left(B_{\tilde{X}_{1}}, B_{\bar{X}_{2}}, 2 t\right) . \tag{9.13}
\end{equation*}
$$

Indeed if $B_{\hat{X}_{1}} \subset \mathcal{U}_{k=1}^{s}\left(u_{k}+2 t B_{\tilde{X}_{2}}\right)$ with $u_{k} \in B_{\bar{X}_{1}}$ for every $k$, then for every $x \in B_{\hat{X}_{1}}$ there is a $k$ such that $x-u_{k} \in 2 t B_{\tilde{X}_{2}}$ and hence $x-u_{k} \in 2\left(B_{\tilde{X}_{1}} \cap t B_{\tilde{X}_{2}}\right)$. We estimate $E\left(B_{\tilde{X}_{1}} \cap t B_{\tilde{X}_{2}}, B_{\tilde{X}_{2}}, t / 2\right)$ by using (4.5) and (4.20). Note that the norm induced by $B_{\tilde{X}_{1}} t B_{\tilde{X}_{2}}$ on $\tilde{X}$ is $\max \left\{\|x\|_{L_{1}(\mu)}, t^{-1}\|x\|_{L_{2}(\mu)}\right\}$. Denote by $\tilde{X}(t)$ the space $\bar{X}$ endowed with this norm. Clearly $d\left(\bar{X}(t), l_{2}^{r}\right) \leqslant t$ and hence by (4.13)

$$
K(\tilde{X}(t)) \leqslant c_{1} \min \left(K(\tilde{X}),(\log t)^{1 / 2}\right)
$$

Also the $\pi_{2}$ norm of the identity map from $\tilde{X}(t)$ to $\tilde{X}_{2}$ is at most that of the identity map from $\tilde{X}_{1}$ into $\tilde{X}_{2}$, i.e. by (4.25), at most $c_{2} n^{1 / 2}$. Thus, we get that

$$
\log E\left(B_{\tilde{X}_{1}} \cap t B_{\tilde{X}_{2}}, B_{\tilde{X}_{2}}, t / 2\right) \leqslant c_{3} n t^{-2} \min \left(K(\tilde{X})^{2}, \log t\right)
$$

Substituting into (9.12) we get that

$$
\begin{equation*}
\log E\left(B_{\tilde{X}_{1}}, B_{\tilde{X}_{2}}, t\right) \leqslant \log E\left(B_{\tilde{X}_{1}}, B_{\tilde{X}_{2}}, 2 t\right)+c_{3} n t^{-2} \min \left(K(\tilde{X})^{2}, \log t\right) \tag{9.14}
\end{equation*}
$$

Assertion (9.11) follows directly from (9.14) by iteration.

We can now easily deduce the main result of this section, which generalizes (9.1) to a general $n$-dimensional subspace of $l_{1}^{N}$ in a 'good position'.

Theorem 9.7. Let $X$ be an n-dimensional subspace of $t_{1}^{N}$. Then there is a probability measure $\mu$ on $\{1, \ldots, N\}$ and a subspace $\tilde{X}$ of $L_{1}(\mu)$ which is isometric to $X$ so that

$$
\begin{equation*}
\log E\left(B_{\tilde{X}_{1}}, B_{\bar{X}_{\infty}}, t\right) \leqslant c n t^{-1}(\log t \cdot \log (N t / n))^{1 / 2}, \quad 2 \leqslant t \leqslant n \tag{9.15}
\end{equation*}
$$

Proof. Choose $\mu$ and $\tilde{X}$ as in Proposition 9.6. Then by (9.10) we may apply Proposition 9.1 to $\tilde{X}$ with $\varrho=2$. By (9.3) and (9.11) we get for $2<s<t$

$$
\begin{aligned}
\log E\left(B_{\tilde{X}_{1}}, B_{\tilde{X}_{\infty}}, t\right) & \leqslant \log E\left(B_{\tilde{X}_{1}}, B_{\tilde{X}_{2}}, t / s\right)+\log E\left(B_{\tilde{X}_{2}}, B_{\tilde{X}_{\alpha}}, s\right) \\
& \leqslant c n\left(t^{-2} s^{2} \log t+s^{-2} \log s N / n\right)
\end{aligned}
$$

By taking $s$ so that $s^{2} \approx t\left(\log (t N / n) \cdot(\log t)^{-1}\right)^{1 / 2},(9.15)$ follows.
Added in proof. It is now known that, up to a possible logarithmic factor, the estimate given in Theorem 6.5 is the best possible. More precisely; there is (for $n \geqslant 2$ ) a constant $C_{n}$ so that for every $0<\varepsilon<1 / 2$ there is a zonotope $P(n, \varepsilon)$ with $B^{n} \subset P(n, \varepsilon) \subset(1+\varepsilon) B^{n}$ and so that the number $N\left(B^{n}, \varepsilon\right)$ of summands of $P(n, \varepsilon)$ satisfies

$$
N\left(B^{n}, \varepsilon\right) \leqslant C_{n}\left(\varepsilon^{-2} \log |\varepsilon|\right)^{(n-1) /(n+2)}
$$

The proof (as well as results for other zonoids) is presented in: J. Bourgain and J. Lindenstrauss, Distributions of points on spheres and approximation by zonotopes, Israel Journal of Mathematics (to appear).

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[^1]:    (1) See also the note "added in proof".

[^2]:    ${ }^{1}$ ) See also the note "added in proof".

