# COMPLEMENTED HILBERTIAN SUBSPACES IN REARRANGEMENT INVARIANT FUNCTION SPACES 

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## Introduction

A classical paper of Rodin and Semenov [RS] studies the closed subspace generated by Rademacher elements in a symmetric space $X$ (defined on the interval $[0,1]$ ) and gives a necessary and sufficient condition on $X$ for this subspace to be isomorphic to the space $l_{2}$. Let $M$ be the Orlicz function defined by $M(u)=e^{u^{2}}-1$, and $L_{M}$ be the associated Orlicz space on $[0,1]$, $\|\quad\|_{M}$ its norm. Then this condition reads:

$$
\begin{equation*}
\exists C<\infty, \forall f \in L_{\infty}([0,1]), \quad\|f\|_{X} \leq C\|f\|_{M} \tag{1}
\end{equation*}
$$

or equivalently the closure $\mathscr{G}$ of $L_{\infty}([0,1])$ in $L_{M}$ is (algebraically) included in $X$.

When this condition is realized, this $l_{2}$ subspace is complemented in $X$ if moreover $X$ is (algebraically) included in the dual Orlicz space $\mathscr{G}^{*}=$ $L_{M_{*}}([0,1])\left(M_{*}\right.$ is the Young conjugate of the function $M$; it is given, up to equivalence, by $M_{*}(t) \simeq t \sqrt{\log e t}$ ), or equivalently:

$$
\begin{equation*}
\exists C<\infty \forall f \in X,\|f\|_{M_{*}} \leq C\|f\|_{X} . \tag{2}
\end{equation*}
$$

This was shown independently by Rodin and Semenov (1979) [RS2] and Lindenstrauss and Tzafriri (1979) [LT2]. These results were extended by Bravermann (1982) [B] in a short note, showing that if a sequence ( $X_{i}$ ) of independent individually distributed random variables spans $l_{2}$ in the rearrangement invariant space $X$, then the variable $X_{i}$ belong to $L_{2}$ and $\mathscr{G} \subset X$. If moreover $\left(X_{i}\right)$ span a complemented closed space, then $X \subset \mathscr{G}^{*}$.

In the first two sections of this paper we show that, roughly speaking, condition (1) characterizes when the r.i. space $X$ contains a subspace isomorphic to $l_{2}$ (in short "hilbertian subspace") while conditions (1) and (2) characterize the situation where $X$ contains a complemented hilbertian subspace. This gives an answer to a question of E. M. Semenov (as reformulated in [T]).

However, as is well known, there exist Orlicz functions $\varphi$ which are not majorized by $M$, (resp. nor majorizing $M_{*}$ ) but such that $L_{\varphi}$ contains a hilbertian (resp. complemented hilbertian) subspace; the elements of the $l_{2}$ basis being disjoint elements of $L_{\varphi}$ (in fact this can be obtained in the sequence space $l_{\varphi}$-see [LT]). Thus we have to discard this case by an additional hypothesis, to be able to obtain (1) and (2) as necessary conditions. We add also an order continuity hypothesis on $X$ (note that if $X$ is not order-continuous, it contains $l_{\infty}$, hence $l_{2}$ ).

The characterization of the existence of hilbertian subspace by condition (1) was already obtained by E. V. Tokarev [T], using a result of Gaposhkin [G] and Rosenthal's embedding theorem for hilbertian subspaces of $L_{1}$. The proof given here is reformulated in order to extend to the case of $l_{p}$-subspaces ( $1<p \leq 2$ ), and to introduce to the proof of the complemented case. This second case was partially solved by Tokarev in the same paper, in fact for subspaces belonging to a particular class $K_{0}(X)$ (in our terminology, when the unit ball of the space is $X$-equiintegrable). In the case where $X$ is $q$-concave, for a $q<2$, this implies the conclusion of our Th. 4, since all hilbertian subspaces of $X$ are then in the class $K_{0}(X)$. It is possible in fact to prove directly (using Lemma 10 below) that, under our hypotheses, if $X$ contains a complemented hilbertian subspace, then it contains another one which belongs to the class $K_{0}(X)$. This gives an alternative proof of our result (Theorem 4), which however does not lead to the quantitative version we give in §3.

In this section, given the additional assumption that $X$ does not contain $c_{0}$, we show the equivalence of the existence of a factorization

$$
l_{2} \xrightarrow{i} X \xrightarrow{\pi} l_{2}
$$

of the identity of the Hilbert space with $\|\pi\|\|i\| \leq C$ and of the existence of two variables $A . G \in X, B . G \in X^{\prime}$, where $G$ is à normal Gaussian variable independent from $(A, B)$, and $\langle A, B\rangle=1,\|A . B\|_{X}\|B \cdot G\|_{X^{\prime}} \leq C$.

Note also that the analogous problem, for finite dimensional spaces $l_{2}^{n}$ was considered in the chapter 9 of [Ka], where a criterion is given for a general Banach lattice not to contain $l_{2}^{n}$ 's uniformly complemented. If one defines the constants $d_{n}(X)$ to be the least constants such that $\sum\left\|f_{i}\right\| \leq d_{n}\left\|\sum f_{i}\right\|$ for all disjoint families of $n$ vectors in $X$, this criterion reads:

$$
\lim \inf d_{n} \cdot(\log n)^{-1 / 2}=0
$$

In the case of r.i. spaces, say on [0,1] it is easy to see that under this condition, $l_{2}$ cannot be a sublattice of $X$, and $X$ does not embed algebraically in $L_{M_{*}}$; i.e., the condition (2) is violated. However this criterion seems to be generally stronger than the negation of condition (2).

In the last two sections, we investigate when a rearrangement invariant function space $X$ is isomorphic to its Hilbert valued extension $X\left(l_{2}\right)$. (This question for spaces with unconditional basis is studied in [KaW]). For an Orlicz space this happens exactly when its Boyd indices are non trivial, i.e., when it is reflexive (Section 4). For a $q$-concave, $q<2$, r.i. space over [0, 1], a necessary and sufficient condition is that the lower Boyd index of the space is non-trivial (strictly greater than 1) (Section 5). Note that, by known results, the non-triviality of both Boyd indices is equivalent (for order-continuous r.i. spaces) to the fact that $X$ has an unconditional basis (in fact, that the Haar basis is unconditional): see [LT2], Theorems 2 c 6 and 2c11. This could perhaps give a way to generalize our results to more general r.i. spaces.

We refer to [LT2] for basic facts about rearrangement invariant function spaces. The definition of r.i. space we consider is that of [LT2]; more precisely the function space $X$ (over $I=[0, a]$ or $=[0, \infty)$ ) is r.i. iff i): for every $x \in X$, its rearrangement $x^{*}$ has same norm, and ii): either simple integrable functions are dense in $X$, or $X$ has Fatou property ([LT2], p. 30). If $(\Omega, \mathscr{A}, \mu)$ is an arbitrary measure space, $X(\Omega, \mathscr{A}, \mu)$ is the space of measurable functions on $\Omega$ whose rearrangements are in $X(I)$.

A reference for ultrapowers and ultraproducts in [H].

## 1. $l_{p}$-subspaces of rearrangement invariant spaces

If $X$ is a Banach lattice, we denote by $X^{\prime}$ the Nakano dual of $X$, i.e., the space of $\sigma$-order continuous elements $x^{*}$ of the dual: $x_{n} \in X, x_{n} \downarrow 0 \Rightarrow$ $\left\langle x_{n}, x^{*}\right\rangle \rightarrow 0$. When $X$ is order continuous then $X^{\prime}=X^{*}$ and $X$ embeds isometrically in $X^{\prime \prime}$.

Proposition 1. Let $X$ be an order continuous rearrangement invariant function space, and $1<p \leq 2$. Assume that $X$ does not contain the space $l_{p}$ as sublattice. Then $X$ contains $l_{p}$ as subspace iff $X^{\prime \prime}$ contains a $p$-stable random variable ( $a$ Gaussian variable when $p=2$ ).

Remark 2. Let $X([0,1])$ be the restriction of the space $X$ to the interval $[0,1]$ (in the case it is defined on $[0, \infty)$ ). Let $\mathscr{f}_{p}$, resp $\mathscr{G}$ be the closure of $L_{\infty}([0,1])$ in the weak $L_{p}$ space $L_{p, \infty}$, resp. in the Orlicz space $L_{M}([0,1])$ $\left(M(u)=e^{u^{2}}-1\right)$. Then $X^{\prime \prime}$ contains a $p$-stable, resp. Gaussian variable iff the space $\mathscr{S}_{p}$, resp. $\mathscr{G}$ is (algebraically) included in $X([0,1])$.

This remark is simply a consequence of the classical estimation of the tail of a $p$-stable random variable $\gamma_{p}$ (see [F]): $\mathbf{P}\left(\left|\gamma_{p}\right|>t\right) \simeq t^{-p}$ for $t \rightarrow \infty$. Hence $L_{p, \infty}([0,1])$ is simply the space of functions which are majorized, up to a rearrangement, by an homothet of $\left|\gamma_{p}\right|+1$. In the Gaussian case, note that the Orlicz space $L_{M}([0,1])$ coincides with the Lorentz space $L_{M, \infty}([0,1])$
(which consists of functions which are majorized by a function equimeasurable with an homothet of $|G|+1, G$ a normal Gaussian variable): see [LT2], Thm. 2b4 and its proof.

Remark 3. When $X$ is order continuous, the Fatou property $\left(X=X^{\prime \prime}\right)$ is equivalent to the fact that $X$ does not contain $c_{0}$ as subspace. In this case, again with the hypothesis that $X$ does not contain $l_{p}$ as sublattice, $X$ contains $l_{p}$ iff it contains a $p$-stable (resp. Gaussian) variable.

Proof of Prop. 1. (a) We prove first the necessity.
Let $E$ be a $l_{p}$ subspace of the r.i. space $X$ (defined for instance on $[0, \infty$ ) equipped with Lebesgue measure $\lambda$ ). Then the norm on $X$ is equivalent to the $L_{1}(U)$ norm on elements of $E$, for some measurable subset $U$ of finite measure. For, if not, there is a sequence $\left(f_{n}\right)_{n}$ in the unit sphere of $E$ such that $\forall a>0, \int_{0}^{a}\left|f_{n}\right| d \lambda \rightarrow 0$ as $n \rightarrow \infty$. Using the order continuity of $X$, we deduce that $\forall k,\left\|\left|f_{k}\right| \wedge\left|f_{n}\right|\right\|_{X} \rightarrow 0$ as $n \rightarrow \infty$; after suitable extraction we obtain a $l_{p}$-basis $\left(f_{n}\right)_{n}$ in $E$ and a disjoint sequence $\left(f_{n}^{\prime}\right)_{n}$ which are equivalent for the norm of $X$, a contradiction.

We use now the following fact, which is a consequence of the paper [DCK] by Dacunha-Castelle and Krivine on subspaces of $L_{1}$ (see also [A]); in the hilbertian case, as noticed in [T], it is also a consequence of Gaposhkin's result on the central limit theorem for sequences of functions ([G], Thm. 1.5.1) and of Rosenthal's theorem on subspaces of $L_{p}$ ([R], Thm. 1). Every infinite dimensional $l_{p}$ subspace of a space $L_{1}(U)$ contains a normalized $l_{p}$ sequence of functions whose distributions are asymptotically conditionally $p$-stable (with same parameter) and conditionally independent. More precisely, there exist a superspace $L_{1}(U \times S, \lambda \otimes \sigma)$ (where $\sigma$ is a probability), a function $Y$ in $L_{1}(U \times S, \lambda \otimes \sigma)$ which has conditional $p$-stable distribution (i.e., $Y(\omega, \cdot)$ is a $p$-stable variable for a.e. $\omega \in U$ ), and a normalized sequence $\left(f_{n}\right)$ in $L_{1}(U)$ which converges "weakly conditionally in distribution" to $Y$, in the sense that for every bounded continuous function $\varphi \in$ $C_{b}(\mathbf{R})$,

$$
\lim _{n \rightarrow \infty} \varphi\left(f_{n}\right)=\int_{S} \varphi(Y(\cdot, s)) d \sigma(s)
$$

where the limit is taken in the $\sigma\left(L_{\infty}, L_{1}\right)$ sense; or, equivalently,

$$
\forall f \in L_{1}(U), \quad\left(f, f_{n}\right) \xrightarrow{\text { dist }}(f, Y)
$$

where "dist" refers to the usual weak convergence of probability distributions (or more generally of finite measures). This second definition is only apparently stronger than the first one, in the case of $L_{1}$-bounded sequences. We
refer to [BR] for an extensive study of this kind of distributional convergence (which we call in short "wcd" hereafter).

We shall use the whole information given by this result only in §2, and use here only the fact that the distribution of $\left(f_{n}\right)$ converges to that of $Y$. Using [LT2], remark following 1.b.18, it is easy to see that $Y \in X^{\prime \prime}$ (with $\|Y\|_{X^{\prime \prime}} \leq$ $\liminf _{n \rightarrow \infty}\left\|f_{n}\right\|_{X}$ ). But $Y$ is equimeasurable with a function of the form $A \otimes \gamma_{p}$, where $A \in L_{1}^{+}(U)$ and $\gamma_{p}$ is a $p$-stable variable (defined on $(S, \sigma)$ ). Since conditional expectation operators act on $X^{\prime \prime}$, we see that

$$
\alpha \mathbf{1}_{U} \otimes \gamma_{p} \in X^{\prime \prime}(U \times S) \quad\left(\text { where } \alpha=\frac{1}{\lambda(U)} \int_{U} A d \lambda\right)
$$

hence $\gamma_{p} \in X^{\prime \prime}(S, \sigma)$.
(b) The sufficiency of the condition results in the hilbertian case from the Rodin-Semenov theorem (see [LT2], Thm. 2b4). In the $l_{p}$-case, we have to work a little bit more. Suppose that $X^{\prime \prime}$ contains a $p$-stable random variable $\gamma_{p}$. It is well known that $\gamma_{p}$ can be realized as a product $G \otimes \gamma^{1 / 2}$, where $G$ is a normal Gaussian variable, and $\gamma$ is a positive ( $p / 2$ )-stable variable. Consider in $X(\Omega \times S)$ the sequence $\left(G_{n} \otimes Y_{n}\right)_{n=1}^{\infty}$, where the $G_{n}$ are independent normal Gaussian variable, and the $Y_{n}$ are independent truncated square roots of ( $p / 2$ )-stable positive variables $\left(Y_{n}^{2}\right.$ is equimeasurable with $\left.\gamma \mathbf{1}_{\{|\gamma| \leq n\}}\right)$. For all $f \in X^{\prime \prime}(\Omega \times S)$, we have

$$
f+G_{n} \otimes Y_{n} \xrightarrow{\text { dist }} f+\Gamma
$$

where $\Gamma$ is a $p$-stable variable independent from $f$ (to fix ideas, let us consider that $\Gamma \in X^{\prime \prime}(T)$ ). Hence

$$
\liminf _{n \rightarrow \infty}\left\|f+G_{n} \otimes Y_{n}\right\|_{X^{\prime \prime}(\Omega \times S)} \geq\|f+\Gamma\|_{X^{\prime \prime}(\Omega \times S \times T)}
$$

But if $f$ takes itself the form $G \otimes h, G$ Gaussian in $L^{0}(\Omega), h \in L^{0}(S)$, then

$$
\begin{aligned}
\left\|G \otimes h+G_{n} \otimes Y_{n}\right\| & =\left\|G \otimes\left(h^{2}+Y_{n}^{2}\right)^{1 / 2}\right\| \leq\left\|G \otimes\left(h^{2}+\gamma\right)^{1 / 2}\right\| \\
& =\left\|G \otimes h+G_{n} \otimes \gamma^{1 / 2}\right\|=\|G \otimes h+\Gamma\|
\end{aligned}
$$

thus we have in fact the equality

$$
\lim _{n \rightarrow \infty}\left\|f+G_{n} \otimes Y_{n}\right\|_{X^{\prime \prime}(\Omega \times S)}=\|f+\Gamma\|_{X^{\prime \prime}(\Omega \times S \times T)}
$$

The end of the reasoning is now a matter of folklore, inspired from [KM] (see also the proof of Lemma 10). Let $E=\overline{\operatorname{span}}\left[G_{n} \otimes Y_{n}\right]_{n \geq 1}$. The preceding shows that the spreading model generated by the sequence $\left(G_{n} \otimes Y_{n}\right)_{n}$ over
$E$ is isometric to the space generated by $E$ and a sequence $\left(\Gamma_{n}\right)_{n}$ of independent $p$-stable random variables (in a suitable extension $X^{\prime \prime}(\Omega \times S \times$ $\left.T^{\mathrm{N}}\right)$ ); i.e.,

$$
\forall x \in E,\left\|x+\sum_{j=1}^{n} a_{j} \Gamma_{j}\right\|=\lim _{k_{n} \rightarrow \infty} \lim _{k_{n-1} \rightarrow \infty} \cdots \lim _{k_{1} \rightarrow \infty}\left\|x+\sum_{j=1}^{n} a_{j} G_{k_{j}} \otimes Y_{k_{j}}\right\|
$$

We have moreover that, for all $k_{1}<k_{2}<\ldots k_{n}$,

$$
\begin{gather*}
\left\|x+\sum_{j=1}^{n} a_{j} G_{k_{j}} \otimes Y_{k_{j}}+a_{n+1} G_{m} \otimes Y_{m}+\sum_{j>n+1} a_{j} \Gamma_{j}\right\|  \tag{3}\\
\underset{m \rightarrow \infty}{ }\left\|x+\sum_{j=1}^{n} a_{j} G_{k_{j}} \otimes Y_{k_{j}}+\sum_{j>n} a_{j} \Gamma_{j}\right\|
\end{gather*}
$$

and this convergence is uniform for $\sum\left|a_{k}\right|^{p} \leq 1$. Note that in (3), the lefthand side is simply

$$
\left\|x+\sum_{j=1}^{n} a_{j} G_{k_{j}} \otimes Y_{k_{j}}+a_{n+1} G_{m} \otimes Y_{m}+\left(\sum_{k>m+1}\left|a_{k}\right|^{p}\right)^{1 / p} \Gamma_{2}\right\|
$$

while the righthand side equals

$$
\left\|x+\sum_{j=1}^{n} a_{j} G_{k_{j}} \otimes Y_{k_{j}}+a_{m+1} \Gamma_{1}+\left(\sum_{k>m+1}\left|a_{k}\right|^{p}\right)^{1 / p} \Gamma_{2}\right\| .
$$

If we choose $n_{k+1}$ such that the difference of the two sides in (3) is less than $\varepsilon 2^{-k}$ (for all ( $a_{k}$ ) with $\sum\left|a_{k}\right|^{p} \leq 1$ ), this procedure gives a subsequence $\left(G_{n_{k}} \otimes Y_{n_{k}}\right)$ which is $(1+\varepsilon) /(1-\varepsilon)$ isomorphic to $l_{p}$.

## 2. Complemented hilbertian subspaces of r.i. spaces

The main result of this section is the following Theorem 4 on complemented embeddability of $l_{2}$ in an r.i. space. In the remainder of this section we associate to such a complemented hilbertian subspace conditionally Gaussian variables in both the spaces $X$ and $X^{\prime}$, which will be the main tool in Sections 3 and 5.

We introduce two notations:
If $X$ is an r.i. space over $[0, \infty)$, we denote by $X([0,1])$ its restriction to the interval $[0,1]$.

If $f, g$ are two functions in $X$, resp. $X^{\prime}$, we denote by $\langle f, g\rangle$ the duality bracket, i.e., the integral $\int f g d \lambda$.

Theorem 4. Let $X$ be an order continuous rearrangement invariant function space (over $[0,1]$ or $[0, \infty)$ ), not containing $l_{2}$ as complemented sublattice. Then $X$ contains a complemented hilbertian subspace iff both $X^{\prime}$ and $X^{\prime \prime}$ contain a Gaussian variable, or equivalently $\mathscr{G} \subset X([0,1]) \subset \mathscr{g}^{\prime}$.

Proof of Theorem 4. The sufficiency of this condition results from the fact that it implies that Rademacher functions span a complemented hilbertian subspace (see [LT2],Thm. 2b4). Now we prove the necessity. By the proof of Proposition 1, it suffices to prove that both $X$ and $X^{\prime}$ must contain an hilbertian subspace which is strongly embedded (i.e., on which the topology of $X$, resp. $X^{\prime}$, agrees with the $L_{1}(A)$-topology, relative to some integrable subset $A$ of $[0, \infty)$ ). Assume that this is not the case. Consider a projection $P$ : $X \rightarrow X$, the range of which is an hilbertian subspace, generated by a sequence $\left(g_{n}\right)_{n=1}^{\infty}$, which is equivalent to the natural $l_{2}$-basis. Since $X^{*}=X^{\prime}$, the projection $P$ takes the form

$$
P f=\sum_{n=1}^{\infty}\left\langle f, h_{n}\right\rangle g_{n}
$$

where $\left(h_{n}\right)_{n=1}^{\infty}$ is a sequence in $X^{\prime}$ which is biorthogonal to the sequence $\left(g_{n}\right)_{n=1}^{\infty}\left(\left\langle g_{n}, h_{m}\right\rangle=\delta_{n m}\right)$ and which is clearly equivalent to the $l_{2}$-basis. Suppose for instance that $X^{\prime}$ does not contain a strongly embedded hilbertian subspace. Then for all $N \geq 1$, there exists some $h \in \overline{\operatorname{span}}\left(h_{k}\right)_{k \geq N}$ such that $\|h\|_{X^{\prime}}=1$ and $\left\|1_{[0, N]} h\right\|_{L_{1}}<2^{-N}$. Proceeding inductively, we obtain a sequence of functions $\bar{h}_{n}=\sum_{j \in J_{n}} \alpha_{j}^{(n)} h_{j}$ in $X^{\prime}$, which are disjoint successive blocks of the $h_{j}$, satisfying $\left\|\bar{h}_{n}\right\|_{X^{\prime}}=1,\left\|\mathbf{1}_{[0, n]} h_{n}\right\|_{L_{1}}<2^{-n}$. We have $\left\|\boldsymbol{\alpha}^{(n)}\right\|_{2}:=\left(\sum_{j \in J_{n}}\left|\alpha_{j}^{(n)}\right|^{2}\right)^{1 / 2} \sim 1$. Choose $\boldsymbol{\beta}^{(n)}:=\left(\beta_{j}^{(n)}\right)_{j} \in l_{2}$ with $\left\|\boldsymbol{\beta}^{(n)}\right\|_{2}=$ $\left\|\boldsymbol{\alpha}^{(n)}\right\|_{2}^{-1}$ and $\sum_{j} \beta_{j}^{(n)} \alpha_{n}^{(n)}=1$, and set $\bar{g}_{n}:=\sum_{j \in J_{n}} \beta_{j}^{(n)} g_{j}$. Then $\bar{P}: X \rightarrow X$ : $\bar{P} f=\Sigma_{n}\left\langle f, \bar{h}_{n}\right\rangle \bar{g}_{n}$ is another projection onto an hilbertian subspace of $X$, with moreover $\bar{h}_{n} \rightarrow 0$ locally in measure as $n \rightarrow \infty$ (i.e., in measure on every integrable subset of $[0, \infty)$ ). Since the unit ball of every r.i. space is bounded in measure, so is the sequence $\left(\bar{g}_{n}\right)_{n}$, and hence $\bar{g}_{n} \bar{h}_{n} \rightarrow 0$ locally in measure as $n \rightarrow \infty$. On the other hand $\int \bar{g}_{n} \bar{h}_{n} d \lambda=\left\langle\bar{g}_{n}, \bar{h}_{n}\right\rangle=1$. The same is true when $X$ is supposed not to contain a strongly embedded hilbertian subspace.

So we can suppose w.l.o.g. that the sequence $\left(g_{n} \cdot h_{n}\right)_{n}$ (which is bounded in $L_{1}$ ) is not $L_{1}$-equiintegrable. By passing if necessary to a subsequence, we can suppose that there exist a $\delta>0$ and disjoint sets $\left(A_{n}\right)$ so that: $\int_{A_{n}}\left|f_{n}\right|\left|g_{n}\right| d \lambda=c_{n}>\delta$. Now, the formula

$$
Q f:=\sum_{n} c_{n}^{-1}\left\langle f, \mathbf{1}_{A_{n}}\right| h_{n}| \rangle \mathbf{1}_{A_{n}}\left|g_{n}\right|
$$

defines a bounded projection in $X$, whose range is a sublattice isomorphic to
$l_{2}$, which provides a contradiction. In fact, choosing unimodular elements $u, v \in L_{\infty}$ with $\mathbf{1}_{A_{n}}\left|h_{n}\right|=\mathbf{1}_{A_{n}} u h_{n}, \mathbf{1}_{A_{n}}\left|g_{n}\right|=\mathbf{1}_{A_{n}} v g_{n}$, we have

$$
Q f=v \sum_{n} c_{n}^{-1} \mathbf{1}_{A_{n}} P\left(\mathbf{1}_{A_{n}} u f\right)=v R(u f)
$$

Up to the coefficients $c_{n}^{-1}$, the operator $R$ is "block-diagonally" extracted from $P$, hence bounded by a well-known argument due to Tonge (see [LT], Prop. 1c8), and so is $Q$. On the other hand the (unconditional) basic sequences $g_{n}^{\prime}:=\mathbf{1}_{A_{n}} g_{n}$ and $h_{n}^{\prime}:=\mathbf{1}_{A_{n}} h_{n}$ are dominated by the $\left(l_{2}\right)$ sequences ( $g_{n}$ ), resp. ( $h_{n}$ ). For, denoting by $\left(\varepsilon_{n}\right)$ a sequence of independent Bernouilli random variables, we have for every sequence $\left(\lambda_{n}\right)$ of scalars (with finite support):

$$
\begin{aligned}
\left\|\sum_{n} \lambda_{n} g_{n}^{\prime}\right\|_{X} & =\left\|\mathbf{E}_{\varepsilon}\left|\sum_{n} \varepsilon_{n} \lambda_{n} g_{n}^{\prime}\right|\right\|_{X} \leq\left\|\mathbf{E}_{\varepsilon}\left|\sum_{n} \varepsilon_{n} \lambda_{n} g_{n}\right|\right\|_{X} \\
& \leq \mathbf{E}_{\varepsilon}\left\|\sum_{n} \varepsilon_{n} \lambda_{n} g_{n}\right\|_{X} \sim\left\|\sum_{n} \lambda_{n} g_{n}\right\|_{X}
\end{aligned}
$$

where the first inequality is a consequence of the pointwise inequality

$$
\mathbf{E}_{\varepsilon}\left|\sum_{n} \varepsilon_{n} \lambda_{n} g_{n}^{\prime}\right| \leq \mathbf{E}_{\varepsilon}\left|\sum_{n} \varepsilon_{n} \lambda_{n} g_{n}\right|
$$

(as long as $\left|g_{n}^{\prime}\right| \leq\left|g_{n}\right|$ pointwise), while the second one is simply the triangular inequality. Thus the sequences $\left(g_{n}^{\prime}\right)$ and $\left(c_{n}^{-1} h_{n}^{\prime}\right)$ are both dominated by the $l_{2}$-basis, and since they are biorthogonal, they are in fact equivalent to the $l_{2}$-basis.

Proposition 5. Let $X$ be an r.i. space over $\Omega=[0,1]$ or $[0, \infty)$, satisfying the conditions of Theorem 4. For every projection $P: X \rightarrow X$ with hilbertian range $E$, there exist an $l_{2}$-basic sequence $\left(x_{n}\right)$ in $E$ and a biorthogonal $l_{2}$-basic sequence $\left(x_{n}^{\prime}\right)$ in the range of $P^{*}$ such that:
(i) the sequence $\left(x_{n}\right)$ converges wcd to a conditionally gaussian variable $A \otimes G \in X^{\prime \prime}(\Omega \times[0,1]) ;$
(ii) the sequence ( $x_{n}^{\prime}$ ) converges wcd to a conditionally gaussian variable $B \otimes G \in X^{\prime}(\Omega \times[0,1]) ;$
(iii) $\langle A, B\rangle>0$.

Proof. The reasoning of the proof of Thm. 4 shows that if $P$ is defined by $P f=\sum_{n=1}^{\infty}\left\langle f, h_{n}\right\rangle g_{n}$, where $\left(g_{n}\right)$ is a $l_{2}$-sequence, there exists an integer $n_{0}$ such that $\left(g_{n}\right)_{n \geq n_{0}}$ and $\left(h_{n}\right)_{n \geq n_{0}}$ span closed spaces of $X$, resp. $X^{\prime}$ whose topology coincides with that of some $L_{1}(A)$ ( $A$ of finite measure). By

Gaposhkin's Theorem there are sequences of successive disjoint $l_{2}$-normalized blocks $x_{n}$ (resp. $\left.B \otimes G\right)$. We have $A \otimes G \in X^{\prime \prime}(\Omega \times[0,1])$ and $B \otimes G$ $\in X^{\prime}(\Omega \times[0,1])$. But in fact Gaposhkin's result is that, after extraction of a subsequence, every system of block coefficients $\boldsymbol{\alpha}^{(n)}:=\left(\alpha_{k}^{(n)}\right)_{k}$ with $\left\|\boldsymbol{\alpha}^{(n)}\right\|_{2}=$ 1 and $\left\|\boldsymbol{\alpha}^{(n)}\right\|_{\infty} \rightarrow 0$ gives rise to this wod convergence. So we may take $l_{2}$-conjugate systems of block-coefficients for the $x_{n}^{\prime}$ 's and the $x_{n}^{\prime}$ 's, and obtain that $\left\langle x_{n}, x_{n}^{\prime}\right\rangle=\delta_{n m}$, i.e., these sequences are biorthogonal.

It is clear that we may suppose that $A, B \geq 0$, so to prove (iii), it suffices to prove that the functions $A$ and $B$ are not disjoint.

Suppose at the contrary that $A$ and $B$ are disjoint, i.e. that there exists a set $U \in \mathscr{A}$ such that $\mathbf{1}_{U^{c}} A=0$ and $\mathbf{1}_{U} B=0$. We then have

$$
\mathbf{1}_{U^{c}} x_{n} \xrightarrow{\mathrm{wcd}} 0 \quad \text { and } \quad \mathbf{1}_{U} x_{n}^{\prime} \xrightarrow{\mathrm{wcd}} 0
$$

a fortiori this convergence happens in distribution, hence in measure. Thus the $L_{1}$-bounded sequences $\mathbf{1}_{U} x_{n} x_{n}^{\prime}$ and $\mathbf{1}_{U^{c}} x_{n} x_{n}^{\prime}$ converge to zero in measure, but at least one of them does not converge to zero in norm (since $\int x_{n} x_{n}^{\prime} d \lambda=1$ ). As in the proof of Thm. 4, we can then exhibit two biorthogonal sequences $\left(y_{n}\right)$ and ( $y_{n}^{\prime}$ ) in $X$ resp. $X^{\prime}$, which are dominated by $\left(x_{n}\right)$, resp. ( $x_{n}^{\prime}$ ), (hence equivalent to the $l_{2}$-basis), and give rise to a projection $X \rightarrow X$ whose range is a complemented hilbertian sublattice of $X$, a contradiction.

Before the end of this section, and in close relation with Prop. 5, we give a result on the projection onto the span of conditionally independent Gaussian variables, which will be used in Sections 3 and 5.

Proposition 6. Let $X$ be an r.i. function space over a product space $(\Omega \times S, \overline{\mathscr{A}} \Sigma, \mu \otimes \sigma)$. Let $G$ be a normal Gaussian variable, and $\left(G_{n}\right)$ a sequence of independent normal Gaussian variables in $L^{0}(S, \Sigma, \sigma)$.
(a) If there are $A, B$ in $L^{0}(\Omega, \mathscr{A}, \mu)$ such that $A \otimes G \in X, B \otimes G \in X^{\prime}$ with $\langle A, B\rangle=1$, then the sequence $\left(A \otimes G_{n}\right)_{n}$ spans in $X$ a $C$-complemented closed space, where $C=\|A \otimes G\|_{X}\|B \otimes G\|_{X^{\prime}}$.
(b) Suppose that $X$ is order-continuous. Then conversely if $A \in L^{0}(S, \Sigma, \sigma)$ is such that the sequence $\left(A \otimes G_{n}\right)_{n=1}^{\infty}$ spans in $X(\Omega \times S)$ a complemented closed subspace, then there exists a function $B \in L^{0}(\Omega, \mathscr{A}, \mu)$ such that $B \otimes G \in X^{\prime}(\Omega \times \Sigma),\langle A, B\rangle=1$, and

$$
\|A \otimes G\|_{X}\|B \otimes G\|_{X^{\prime}} \leq C
$$

In particular, $X^{\prime}$ contains a Gaussian variable.
Proof. (a) The sequences $\left(A \otimes G_{n}\right)$ and ( $B \otimes G_{n}$ ) span isometric copies of $l_{2}$ in $X$, resp $X^{\prime}$. We define a projection $R: X \rightarrow X$ by: $R f=\sum_{n=1}^{\infty}\langle f$,
$B \otimes G_{n}>A \otimes G_{n}$. The norm of $R$ is evaluated as follows:

$$
\begin{aligned}
\|R f\|_{X} & =\left\|\sum_{n}\left\langle f, B \otimes G_{n}\right\rangle A \otimes G_{n}\right\|_{X} \\
& =\left(\sum_{n}\left|\left\langle f, B \otimes G_{n}\right\rangle\right|^{2}\right)^{1 / 2}\|A \otimes G\|_{X} \\
& =\left(\sum_{n} \alpha_{n}\left\langle f, B \otimes G_{n}\right\rangle\right)\|A \otimes G\|_{X} \quad \text { for some } \alpha_{n} \in \mathbf{R}, \sum_{n} \alpha_{n}^{2}=1 \\
& =\left\langle f, \sum_{n} \alpha_{n} B \otimes G_{n}\right\rangle\|A \otimes G\|_{X} \\
& \leq\|f\|_{X}\left\|\sum_{n} \alpha_{n} B \otimes G_{n}\right\|_{X^{\prime}}\|A \otimes G\|_{X} \\
& =\|f\|_{X}\|B \otimes G\|_{X^{\prime}}\|A \otimes G\|_{X}
\end{aligned}
$$

Hence $\|R f\| \leq C\|f\|$.
(b) Now we prove the converse. We can w.l.o.g. suppose that $(S, \Sigma, \mu)$ is the product space $\mathbf{R}^{\mathbf{N}}$ equipped with the standard Gaussian measure $\gamma$, and that $G_{n}$ is the $n$th coordinate map $\mathbf{R}^{\mathbf{N}} \rightarrow \mathbf{R}$. The orthogonal group $O(n)$ acts on $\mathbf{R}^{n}$, leaving the $n$-dimensional gaussian measure invariant; let us consider that $O(n)$ acts on $\mathbf{R}^{\mathbf{N}}$, by changing only the $n$ first coordinates. Each element $U$ of $O(n)$ gives rise to an isometry of $X$, again denoted by $U$, and defined by

$$
U . f\left(\omega,\left(x_{1}, \ldots, x_{n}\right), x_{n+1}, \ldots\right)=f\left(\omega, U^{*}\left(x_{1}, \ldots x_{n}\right), x_{n+1} \ldots\right)
$$

Note that if $\left(u_{i j}\right)$ is the matrix of $U$ (relatively to the natural basis) and $f=\sum_{j=1}^{n} \lambda_{j} \otimes G_{j}$ (where the $\lambda_{i}$ are $\mathscr{A}$-measurable functions) then $U . f=$ $\sum_{i}\left(\sum_{j} u_{i j} \lambda_{j}\right) \otimes G_{i}$.

Let $E=\overline{\operatorname{span}}\left[A \otimes G_{n}\right]_{n=1}^{\infty}$ and $P$ be a given projection from $X$ onto $E$, and set

$$
R_{n}=\int_{O(n)} U^{*} P U d \sigma_{n}(U)
$$

where $\sigma_{n}$ is the normalized Haar measure on the compact group $O(n) . \quad R_{n}$ is clearly a projection onto $E$, invariant under the action of $O(n)$ (i.e., $\left.\forall V \in O(n), R_{n} V=V R_{n}\right)$ and of norm $\left\|R_{n}\right\| \leq\|P\|$. Note that $R_{n} R_{m}=$ $R_{m} R_{n}=R_{m \wedge n}$.

The set $\left(R_{n}\right)_{n=1}^{\infty}$ is relatively compact in the weak operator topology (due to the reflexivity of $E$ ), and has a unique cluster point $R$ (because if $f \in X$
depends only on $\omega$ and the $n$ first coordinates, then $\forall m \geq n, R_{m} f=R_{n} f$ ), which is invariant under the action of $O(n)$ for all $n \in \mathbf{N}$. Let us write:

$$
R f=\sum_{i=1}^{\infty}\left\langle f, h_{i}\right\rangle A \otimes G_{i}
$$

where $h_{i} \in X^{*}\left(\Omega \times \mathbf{R}^{\mathbf{N}}\right)$ (note that $X^{\prime}=X^{*}$ as $X$ is supposed to be order continuous). Then for every $U \in O(n)$,

$$
\begin{aligned}
R U f & =\sum_{i=1}^{\infty}\left\langle U f, h_{i}\right\rangle A \otimes G_{i} \\
& =\sum_{i=1}^{\infty}\left\langle f, U^{*} h_{i}\right\rangle A \otimes G_{i}=\sum_{i=1}^{\infty}\left\langle f, h_{i} \circ U\right\rangle A \otimes G_{i}
\end{aligned}
$$

and, on the other hand, $\left(u_{i j}\right)$ being the matrix of $U$,

$$
U R f=\sum_{i=1}^{\infty}\left(\sum_{j} u_{i j}\left\langle f, h_{j}\right\rangle A \otimes G_{j}\right)=\sum_{i=1}^{\infty}\left\langle f, \sum_{j} u_{i j} h_{j}\right\rangle A \otimes G_{j}
$$

whence we obtain

$$
h_{i} \circ U=\sum_{j=1}^{n} u_{i j} h_{j} .
$$

Considering the sequence $\left(h_{i}\right)_{i=1}^{\infty}$ as a measurable map $h: \Omega \times \mathbf{R}^{\mathbf{N}} \rightarrow \mathbf{R}^{\mathbf{N}}$, and associating to $U \in O(n)$ the bijection $\tilde{U}$ of $\Omega \times \mathbf{R}^{\mathbf{N}}$ acting only on the $n$ first coordinateO]s in $\mathbf{R}^{\mathbf{N}}$ as $U$, we have thus the functional equation: $h \circ \tilde{U}=U \circ h$.

If $U$ belongs to the subgroup $\Gamma_{1, n}$ of $O(n)$, whose elements leave the first coordinate inchanged, we obtain in particular: $h_{1}=h_{1} \circ \tilde{U}$, hence $h_{1}=$ $\int_{\Gamma_{1, n}} h_{1} \circ \tilde{U} d \sigma_{1, n}(U)$, where $\sigma_{1, n}$ is the Haar measure on $\Gamma_{1, n}$, hence $h_{1}$ is clearly a (measurable) function of $\omega, x_{1}, \sum_{j=2}^{n} x_{j}^{2}, x_{n+1}, \ldots$. Let $\mathscr{F}_{n}$ be the measure-complete $\sigma$-algebra generated by $\mathscr{A}, G_{1}, \sum_{j=2}^{n} G_{j}^{2}, G_{n+1} \ldots$ and $\mathscr{F}_{\infty}=\bigcap_{n} \mathscr{F}_{n}$. Then $h_{1}$ is $\mathscr{F}_{n}$-measurable for all $n$, hence $\mathscr{F}_{\infty}$-measurable. But $\mathscr{F}_{\infty}=\bar{\sigma}\left(\mathscr{A}, G_{1}\right)$ (the measure-complete $\sigma$-algebra generated by $\mathscr{A}$ and $G_{1}$ ): for, if $f \in L_{1}$ is of the form $f=g \cdot \varphi\left(G_{2}, \ldots G_{n}\right)$, where $g \in L_{1}$ is $\sigma\left(\mathscr{A}, G_{1}\right)$-measurable, and $\varphi$ is a continuous bounded function on $\mathbf{R}^{n}$, set $f_{k}=g \cdot \varphi\left(G_{2+k n}, \ldots G_{(k+1) n}\right)$; then $\mathbf{E}^{\mathscr{F}_{\infty}} f_{k}=\mathbf{E}^{\mathscr{F}_{\infty}} f$ by symmetry; hence

$$
\mathbf{E}^{\mathscr{F _ { \infty }}} f=\mathbf{E}^{\mathscr{\mathscr { F } _ { \infty }}} \frac{f_{0}+\cdots+f_{k}}{k} \underset{k \rightarrow \infty}{\longrightarrow} g \mathbf{E} \varphi\left(G_{2}, \ldots G_{n}\right)
$$

by the law of large numbers applied to the $f_{k}=\varphi\left(G_{2+k n}, \ldots G_{(k+1) n}\right)$.

Finally $h_{1}\left(\omega, x_{1}, \ldots\right)=H\left(\omega, x_{1}\right)$ a.e., and, using the relation $h \circ \tilde{U}=U \circ h$ for $U$ the orthogonal symmetry exchanging the coordinate $x_{1}$ and $x_{k}$, we have

$$
h_{k}\left(\omega, x_{1}, \ldots x_{k} \ldots\right)=H\left(\omega, x_{k}\right) \text { a.e. }(\forall k \geq 2)
$$

Using now the relation $h \circ \tilde{U}=U \circ h$ for $U$ being the central symmetry, we obtain

$$
\begin{equation*}
H(\omega,-x)=-H(\omega, x) \text { for a.e. }(\omega, x) \tag{1}
\end{equation*}
$$

and if $U$ is the transformation

$$
\left(x_{1}, x_{2}\right) \mapsto\left(\frac{x_{1}+x_{2}}{\sqrt{2}}, \frac{x_{1}-x_{2}}{\sqrt{2}}\right)
$$

we obtain

$$
\begin{equation*}
H\left(\omega, \frac{x_{1} \pm x_{2}}{\sqrt{2}}\right)=\frac{H\left(\omega, x_{1}\right) \pm H\left(\omega, x_{2}\right)}{\sqrt{2}} \quad \text { for a.e. } \omega, x_{1}, x_{2} \tag{2}
\end{equation*}
$$

Then $H$ is a.e. equivalent to a function $\tilde{H}$, whose partial functions $\tilde{H}_{\omega}$ : $x \mapsto \tilde{H}(\omega, x)$ are of class $C^{1}$ for a.e. $\omega$. (if $\varphi$ is a centrally symmetric $C^{1}$ function with compact support and integral 1, it is straightforward that $H_{\omega}(x)=\sqrt{2} H_{\omega} * \varphi(x / \sqrt{2})$ for a.e. $(\omega, x)$. Then, for a.e. $\omega, \tilde{H}$ verifies (1) and (2) for all $x, x_{1}, x_{2}$, and by a standard reasoning, $\tilde{H}(\omega, x)=B(\omega) \cdot x$.

Coming back to the projection $R$, we see that it can be written as

$$
R f=\sum_{i=1}^{\infty}\left\langle g, B \otimes G_{i}\right\rangle A \otimes G_{i}
$$

Note that $A \in X, B \in X^{\prime}$, with $\|A\|_{X} \leq\|A \otimes G\|_{X} / \mathbf{E}|G|$ and $\|B\|_{X^{\prime}} \leq \| B$ $\otimes G \|_{X^{\prime}} / \mathbf{E}|G|$. From $R\left(A \otimes G_{i}\right)=A \otimes G_{i}$ we obtain $\langle A, B\rangle=1$.

## 3. Quantitative version of the preceding results

We say that the space $l_{2}$ is $C$-representable as complemented subspace of $X$ if the identity map of $l_{2}$ is $C$-factorizable through $X$, i.e., there exist linear operators $i: l_{2} \rightarrow X$ and $\pi: X \rightarrow l_{2}$, such that $\pi \circ i=\operatorname{id}_{l_{2}}$ and $\|\pi\|\|i\| \leq C$.

In this section we prove the following improvement of Theorem 4 (with a slight reinforcement of the hypotheses on the r.i. space):

Theorem 7. Let $X$ be a rearrangement invariant function space, not containing $c_{0}$ as subspace nor $l_{2}$ as complemented sublattice. Then $l_{2}$ is $C$-representable in $X$ as complemented subspace iff there exist variables $A \otimes G$,
resp. $B \otimes G$ in $X(\Omega \times[0,1])$, resp. $X^{\prime}(\Omega \times[0,1])$ where $G$ is a normalized Gaussian variable, $\|A \otimes G\|_{X}\|B \otimes G\|_{X^{\prime}} \leq C$ and $\langle A, B\rangle=1$.

The proof of Theorem 7 involves several lemmas. We give first some preliminary material.

A sequence $\left(x_{k}\right)_{k}$ in $X$ is $X$-equiintegrable iff it satisfies the conditions

$$
\begin{equation*}
\lim _{M \rightarrow \infty} \sup _{k}\left\|x_{k} \mathbf{1}_{\left\{\left|x_{k}\right|>M\right\}}\right\|_{X}=0 ; \quad \inf _{\mu(A)<\infty} \sup _{k}\left\|\mathbf{1}_{A^{c}} x_{k}\right\|_{X}=0 \tag{4}
\end{equation*}
$$

As is well known (see for instance [W]) in an ri. space not containing $c_{0}$, every sequence $\left(x_{n}\right)$ has a subsequence $x_{n_{k}}$ for which there is a splitting $x_{n_{k}}=x_{k}^{\prime}+x_{k}^{\prime \prime}$, where each $x_{k}^{\prime}$ is disjoint from the corresponding $x_{k}^{\prime \prime}$, the elements $x_{k}^{\prime \prime}$ are disjoint and the sequence $\left(x_{k}^{\prime}\right)_{k}$ is $X$-equiintegrable.

We shall also use repeatedly in the subsequent proofs the following well known fact, which we will call the "Bessaga-Pelczynski perturbation principle": if $\left(x_{n}\right)$ is a basic unconditional sequence which spans a complemented subspace in $X$, and $\left(y_{n}\right)$ is a sequence in $X$ such that $\left\|x_{n}-y_{n}\right\| \rightarrow 0$, then a subsequence $\left(y_{n_{k}}\right)$ is equivalent to $\left(x_{n_{k}}\right)$, and spans a subspace which is complemented in $X$. In fact, if $P: X \rightarrow \overline{\operatorname{span}}\left[x_{n}\right]$ is a given projection, and $\pi$ : $\overline{\operatorname{span}}\left[x_{n}\right] \rightarrow \overline{\operatorname{span}}\left[x_{n_{k}}\right]$ is the natural projection, then, if $\left(n_{k}\right)$ is sufficiently lacunary, the restriction $J$ of $\pi P$ to $F=\overline{\operatorname{span}}\left[y_{n_{k}}\right]$ is an isomorphism onto $E=\overline{\operatorname{span}}\left[x_{n_{k}}\right]$, and $Q=J^{-1} \pi P$ is a projection. Note that this construction is of almost isometric nature, i.e. we can obtain that the tails $\left(y_{n_{k}}\right)_{k \geq m}$ are $\left(1+\varepsilon_{m}\right)$-equivalent to $\left(x_{n_{k}}\right)_{k \geq m}$ (with $\left.\varepsilon_{m} \rightarrow 0\right)$ and, if $\left(x_{n}\right)$ is $K$-suppression unconditional, the projection $Q_{m}$ onto $F_{m}=\overline{\operatorname{span}}\left[y_{n_{k}}\right]_{k \geq m}$ to be of norm $\leq\left(1+\varepsilon_{m}\right)\|P\| K$.

In the same spirit we state the following very elementary fact, in order to avoid further repetitions:

Lemma 8. Let $F$ be a complemented hilbertian subspace of $X, P: X \rightarrow F a$ projection, $S: F \rightarrow X$ a bounded operator. If PS: $F \rightarrow F$ is an (into) isomorphism, then $S F$ is itself a complemented hilbertian subspace of $X$ (and $S$ an isomorphism).

Proof. $\left.S\right|_{F}$ is an isomorphism, so $G=S F$ is hilbertian; $J=\left.P\right|_{G}$ is an isomorphism from $G$ into $F$. Let $\pi$ be a projection $F \rightarrow P(G)$. Then $Q=J^{-1} \pi P$ is a projection $X \rightarrow G$.

The following lemma precises Proposition 5 in the case $X \not \supset c_{0}$, and will be given a quantitative version by Lemma 12:

Lemma 9. Let $X$ satisfy the hypotheses of Theorem 7. Then for every complemented hilbertian subspace $E$ of $X$ there is a sequence which is $X$-equiin-
tegrable, converges wcd to a conditionally Gaussian variable and spans a complemented hilbertian subspace, and is arbitrarily close in measure to $E$.

Proof. Let $\left(x_{n}\right)_{n}$ be an $l_{2}$-basis of $E$, and $P$ be the projection onto $E$.
First we may suppose that the elements of $E$ live on a $\mu$-finite subset $U$ of $\Omega$ : if there is no $\mu$-integrable subset $U$ such that the norm of $X(\Omega)$ and that of $X(U)$ are equivalent on $U$, then using the order continuity of $X$ we construct recursively a sequence $y_{n}$ of disjoint blocks on the basis $x_{n}$ and a sequence of disjoint elements $y_{n}^{\prime}$ of $X$ such that $\left\|y_{n}-y_{n}^{\prime}\right\| \rightarrow 0$; by the Bessaga-Pelczynski perturbation principle, a subsequence of ( $y_{n}^{\prime}$ ) spans a complemented hilbertian space, a contradiction. Then, reasoning as in §1, we may suppose that the norms of $X$ and of $L_{1}(U)$ are equivalent on $E$, and that the $x_{n}$ converge (weakly conditionally in distribution) to a conditionally Gaussian variable.

By passing if necessary to a subsequence we have a disjoint splitting $x_{k}=x_{k}^{\prime}+x_{k}^{\prime \prime}$, where the elements $x_{k}^{\prime \prime}$ are disjoint and the sequence $\left(x_{k}^{\prime}\right)_{k}$ is $X$-equiintegrable.

Let $E^{\prime}$, resp. $E^{\prime \prime}$ be the closed subspaces of $X$ spanned by the sequence $\left(x_{k}^{\prime}\right)_{k}$, resp. $\left(x_{k}^{\prime \prime}\right)_{k}$. If $x=\sum_{k=1}^{n} \alpha_{k} x_{k}$, set $S^{\prime} x=\sum_{k} \alpha_{k} x_{k}^{\prime}$ and $S^{\prime \prime} x=\sum_{k} \alpha_{k} x_{k}^{\prime \prime}$; note that $x=S^{\prime} x+S^{\prime \prime} x$, and that $S^{\prime}$ and $S^{\prime \prime}$ extend to bounded operators from $E$ to $E^{\prime}$, resp. $E$ to $E^{\prime \prime}$, since the disjoint basic sequence $\left(x_{n}^{\prime \prime}\right)$ is dominated by $\left(x_{n}\right)$; i.e., $\left\|\sum_{k} \alpha_{k} x_{k}^{\prime \prime}\right\| \leq C\left\|\sum_{k} \alpha_{k} x_{k}\right\|$ where $C$ is the unconditionality constant of $\left(x_{n}\right)_{n}$ (see the proof of Thm. 4). Then there exists $M$ and a subspace $E_{1}$ of $E$ of finite codimension (spanned by the $x_{k}, k \geq N_{1}$ ) such that $\forall x \in E_{1},\|x\| \leq M\left\|P S^{\prime} x\right\|$. For if not, there is a sequence $\left(y_{n}\right)_{n}$ in $E$, with $\left\|y_{n}\right\|=1,\left\|P S^{\prime} y_{n}\right\| \rightarrow 0$ and $y_{n} \rightarrow 0$ weakly. Note that $y_{n}-P S^{\prime \prime} y_{n}=$ $P y_{n}-P S^{\prime \prime} y_{n}=P S^{\prime} y_{n} \rightarrow 0$. Again, using the Bessaga-Pelczynski perturbation principle, we find a subsequence $\left(z_{n}\right)_{n}$ of $\left(y_{n}\right)$, such that the sequence $\left(P S^{\prime \prime} z_{n}\right)_{n}$ is basic and equivalent to $\left(z_{n}\right)$. In particular $\left\|P S^{\prime \prime} z\right\| \geq \delta\|z\|$ for all $z \in \operatorname{span}\left[z_{n}\right]$. Since $\overline{\operatorname{span}}\left[z_{n}\right]$ is complemented in $E$, we deduce by Lemma 5 that $\overline{\operatorname{span}}\left[S^{\prime \prime} z_{n}\right]$ is complemented in $X$, and ( $S^{\prime \prime} z_{n}$ ) is equivalent to the $l_{2}$-basis. Hence $X$ contains $l_{2}$ as complemented sublattice, a contradiction.

Thus we find $E_{1}$ on which $\|y\| \simeq\left\|P y^{\prime}\right\|$, hence $\simeq\left\|y^{\prime}\right\|$. Let $E_{1}^{\prime}:=\left\{y^{\prime} / y \in\right.$ $\left.E_{1}\right\}$. Then $E_{1}$ is hilbertian, and moreover is complemented by Lemma 8. Note also that we have

$$
x_{n}^{\prime \prime} \xrightarrow[n \rightarrow \infty]{\text { dist }} 0
$$

So $\left(x_{n}^{\prime \prime}\right)_{n}$ converges to zero in measure, and $\left(x_{n}^{\prime}\right)_{n}$ converges wcd to a conditionally Gaussian variable (the same as for $\left.\left(x_{n}\right)_{n}\right)$.

The key for obtaining a quantitative version of Lemma 9 is the following:
Lemma 10. Let $X$ be an order continuous r.i. function space and $\left(x_{n}\right)_{n}$ a $X$-equiintegrable sequence which converges wcd to a non-zero conditionally
gaussian variable. Then for every $\varepsilon>0$ there exists a subsequence $\left(x_{n_{i}}\right)_{i}$ which is $(1+\varepsilon)$ equivalent to the $l_{2}$ basis and such that the unit ball of $F=\frac{\sin }{\operatorname{span}}\left[x_{n_{i}}\right]_{i}$ is $X$-equiintegrable.

Proof. Choose a sequence of reals $\varepsilon_{n}>0$ with $\sum_{n=1}^{\infty} \varepsilon_{n} \leq \varepsilon$.
The sequence $\left(x_{n}\right)$ converges wcd to $A \otimes G$, where $G \in L_{0}([0,1])$ is a normal gaussian variable and $A \in L_{0}(\Omega)$ is such that $A \otimes G \in X(\Omega \times[0,1])$. We may suppose that $L_{0}([0,1])$ contains an auxiliary normal Gaussian variable $G^{\prime}$, independent from $G$.

Throughout the proof of this lemma, if $x \in X$, we denote by $x^{*}$ the non-increasing rearrangement of $|x|$; this is an element of $X([0, \infty))$ or $X([0,1])$.

Now we construct a subsequence $\left(x_{n_{i}}\right)_{i \geq 1}$ of $\left(x_{n}\right)_{n}$ such that for all $k$ we have the following property, denoted by $\left(H_{k}\right)$ :

$$
\begin{aligned}
& \left\|\left(\sum_{i=1}^{k} \lambda_{i} x_{n_{i}}+\rho A \otimes G\right)^{*}-\left(\sum_{I=1}^{k} \lambda_{i}^{2}+\rho^{2}\right)^{1 / 2}(A \otimes G)^{*}\right\|_{X} \\
& \quad \leq\left(\sum_{i=1}^{k} \varepsilon_{i}\right)\left(\sum_{i=1}^{k} \lambda_{i}^{2}+\rho^{2}\right)^{1 / 2}\|A \otimes G\|_{X} \quad \forall \lambda_{1}, \ldots \lambda_{k}, \rho \in \mathbf{R} .
\end{aligned}
$$

Suppose that we have found the first $k$ terms $n_{1}<n_{2}<\ldots n_{k}$ (possibly $k=0$ ). For every $\lambda_{1}, \ldots \lambda_{k}, \rho \in \mathbf{R}$ we have

$$
\sum_{i=1}^{k} \lambda_{i} x_{n_{i}}+\rho A \otimes G^{\prime}+\lambda_{k+1} x_{n} \xrightarrow[n \rightarrow \infty]{\text { dist }} \sum_{i=1}^{k} \lambda_{i} x_{n_{i}}+\rho A \otimes G^{\prime}+\lambda_{k+1} A \otimes G
$$

which has the same distribution as

$$
\sum_{i=1}^{k} \lambda_{i} x_{n_{i}}+\left(\lambda_{k+1}^{2}+\rho^{2}\right)^{1 / 2} A \otimes G
$$

This implies the convergence Lebesgue-a.e. of the non-increasing rearrangements:

$$
\begin{aligned}
&\left(\sum_{i=1}^{k} \lambda_{i} x_{n_{i}}+\rho A \otimes G^{\prime}+\lambda_{k+1} x_{n}\right)^{*} \\
& \xrightarrow{\text { a.e. }}\left(\sum_{i=1}^{k} \lambda_{i} x_{n_{i}}+\left(\lambda_{k+1}^{2}+\rho^{2}\right)^{1 / 2} A \otimes G\right)^{*}
\end{aligned}
$$

As the sequence $\left(x_{n}\right)$ is $X$-equiintegrable, we deduce that this convergence holds also in the sense of the $X$-norm. Now let

$$
\begin{aligned}
& F_{n}\left(\lambda_{1}, \ldots \lambda_{k+1}, \rho\right)=\left(\sum_{i=1}^{k} \lambda_{i} x_{n_{i}}+\rho A \otimes G^{\prime}+\lambda_{k+1} x_{n}\right)^{*} \\
& F_{\infty}\left(\lambda_{1}, \ldots \lambda_{k+1}, \rho\right)=\left(\sum_{i=1}^{k} \lambda_{i} x_{n_{i}}+\left(\lambda_{k+1}^{2}+\rho^{2}\right)^{1 / 2} A \otimes G\right)^{*} .
\end{aligned}
$$

The $X([0, \infty))$ valued functions $F_{n}$ are equicontinuous on the compact set

$$
K_{k+2}=\left\{\left(\sum_{i=1}^{k+1} \lambda_{i}^{2}+\rho^{2}\right)^{1 / 2} \leq 1\right\} \subset \mathbf{R}^{k+2}
$$

this is a straightforward consequence of the Lipschitz inequality of LorentzShimogaki [LS],

$$
\forall u, v \in X, \quad\left\|u^{*}-v^{*}\right\|_{X} \leq\|u-v\|_{X}
$$

which is true for every r.i. space. Thus the convergence $F_{n} \rightarrow F_{\infty}$ holds uniformly on $K_{k+2}$. We choose $n_{k+1}>n_{k}$ such that $\left\|F_{n_{k+1}}-F_{\infty}\right\|_{\infty} \leq$ $\varepsilon_{k+1}\|A \otimes G\|_{X}$, and we obtain

$$
\begin{aligned}
& \|\left(\sum_{i=1}^{k} \lambda_{i} x_{n_{i}}+\rho A \otimes G^{\prime}+\lambda_{k+1} x_{n}\right)^{*} \\
& -\left(\sum_{i=1}^{k} \lambda_{i} x_{n_{i}}+\left(\lambda_{k+1}^{2}+\rho^{2}\right)^{1 / 2} A \otimes G\right)^{*} \|_{X} \\
& \quad \leq \varepsilon_{k+1}\left(\sum_{i=1}^{k} \lambda_{i}^{2}+\rho^{2}\right)^{1 / 2}\|A \otimes G\|_{X}
\end{aligned}
$$

This inequality together with $\left(H_{k}\right)$ implies $\left(H_{k+1}\right)$. Now if $\Sigma_{i} \lambda_{i}^{2}<\infty$, we apply ( $H_{k}$ ) with $\rho=0$, pass to the limit on $k$, and obtain finally

$$
\left\|\left(\sum_{i=1}^{\infty} \lambda_{i} x_{n_{i}}\right)^{*}-\left(\sum_{i=1}^{\infty} \lambda_{i}^{2}\right)^{1 / 2}(A \otimes G)^{*}\right\|_{X} \leq \varepsilon\left(\sum_{i=1}^{\infty} \lambda_{i}^{2}\right)^{1 / 2}\|A \otimes G\|_{X}
$$

This implies that the sequence $\left(x_{n_{i}}\right)$ is $(1+\varepsilon)$-equivalent to the $l_{2}$ basis.

Moreover, the proof of this inequality gives in fact that

$$
\left\|\left(\sum_{i=k}^{\infty} \lambda_{i} x_{n_{i}}\right)^{*}-\left(\sum_{i=k}^{\infty} \lambda_{i}^{2}\right)^{1 / 2}(A \otimes G)^{*}\right\|_{X} \leq \delta_{k}\left(\sum_{i=k}^{\infty} \lambda_{i}^{2}\right)^{1 / 2}\|A \otimes G\|_{X}
$$

where $\delta_{k}=\sum_{i=k}^{\infty} \varepsilon_{i}$. This implies that, if $\left(\sum_{i=1}^{\infty} \lambda_{i}^{2}\right)^{1 / 2} \leq 1$,

$$
\begin{aligned}
\left\|\mathbf{1}_{U}\left(\sum_{i=1}^{\infty} \lambda_{i} x_{n_{i}}\right)\right\|_{X} \leq & \left\|\mathbf{1}_{U}\left(\sum_{i=1}^{k} \lambda_{i} x_{n_{i}}\right)\right\|_{X}+\delta_{k+1}\|A \otimes G\|_{X} \\
& +\left\|\mathbf{1}_{[0, \mu(U)]}(A \otimes G)^{*}\right\|_{X} \quad \forall U \subset \Omega
\end{aligned}
$$

In the case where $\mu(\Omega)=1$, this implies that the unit ball of $F=\overline{\operatorname{span}}\left[x_{n_{i}}\right]$ is $X$-equiintegrable. In the case $\mu(\Omega)=\infty$, we may suppose that $\mu$ is $\sigma$-finite; let $\left(\Omega_{p}\right)_{p}$ be an increasing sequence of $\mu$-integrable subsets whose union is $\Omega$. Note that each sequence $\left(\mathbf{1}_{\Omega_{p}^{c}} x_{n}\right)_{n}$ is $X$-equiintegrable and wcd converging to $\mathbf{1}_{\Omega_{\rho}^{c}} A \otimes G$. Then a diagonal argument allows us to obtain a subsequence $\left(x_{n_{i}}\right)$ such that for every $k, p$ with $k \geq p$,

$$
\left\|\left(\sum_{i=k}^{\infty} \lambda_{i} \mathbf{1}_{\Omega_{p}^{c}} x_{n_{i}}\right)^{*}-\left(\sum_{i=k}^{\infty} \lambda_{i}^{2}\right)^{1 / 2}\left(\mathbf{1}_{\Omega_{p}^{c}} A \otimes G\right)^{*}\right\|_{X} \leq \delta_{k}\left(\sum_{i=k}^{\infty} \lambda_{i}^{2}\right)^{1 / 2}\left\|\mathbf{1}_{\Omega_{p}^{c}} A \otimes G\right\|_{X}
$$

Hence if $\left(\sum_{i=1}^{\infty} \lambda_{i}^{2}\right)^{1 / 2} \leq 1$ and $q \geq p$ :

$$
\begin{aligned}
\left\|\mathbf{1}_{\Omega_{q}^{c}} \sum_{i=1}^{\infty} \lambda_{i} x_{n_{i}}\right\| & \leq \sum_{i=1}^{p}\left|\lambda_{i}\right|\left\|\mathbf{1}_{\Omega_{q}^{c}} x_{n_{i}}\right\|+\left\|\mathbf{1}_{\Omega_{p}^{c}} \sum_{i=p}^{\infty} \lambda_{i} x_{n_{i}}\right\| \\
& \leq \sqrt{p} \bigvee_{i=1}^{p}\left\|\mathbf{1}_{\Omega_{q}^{c}} x_{n_{i}}\right\|+(1+\varepsilon)\left\|\mathbf{1}_{\Omega_{p}^{c}} A \otimes G\right\|_{X}
\end{aligned}
$$

Letting $q \rightarrow \infty$ and then $p \rightarrow \infty$ we obtain the second condition in (4) for the $X$-equiintegrability of the sequence $\left(x_{n_{i}}\right)$.

Remark 11. We can choose the subsequence ( $x_{n_{i}}$ ) such that every normalized weakly null sequence $\left(z_{l}\right)$ in $\overline{\operatorname{span}}\left[x_{n_{t}}\right]$ converges wad to the same conditionally Gaussian variable.

It is sufficient to prove this fact when $\left(z_{l}\right)$ is a sequence of successive $l_{2}$-normalized blocks on the $x_{n_{i}}$ (since every weakly null sequence $\left(z_{l}\right)$ has a subsequence ( $z_{l_{i}}$ ) which can be approximated in $X$-norm, and a fortiori weakly conditionally in distribution, by such a sequence of disjoint successive blocks on the basis $\left(x_{n}\right)$ ).

It is clear by the preceding proof that for each $V \in \mathscr{A}$, we can choose the subsequence ( $x_{n_{i}}$ ) such that for all $k \geq k_{V}$, for every block $z$ built on the $x_{n_{i}}$, $i \geq k$, we have

$$
\left\|\left(\mathbf{1}_{V} z\right)^{*}-\right\| z\left\|_{2}\left(\mathbf{1}_{V} A \otimes G\right)^{*}\right\| \leq \delta_{k}\|z\|_{2}\left\|\mathbf{1}_{V} A \otimes G\right\|
$$

By a diagonal argument, this can be done for all $V$ in a countable subset $\Gamma$ of $\mathscr{A}$. If the measure space $(\Omega, \mathscr{A}, \mu)$ is separable, this shows that for every $V \in X$, and every sequence of successive disjoint $l_{2}$-normalized blocks $\left(z_{l}\right)$ on the ( $x_{n_{i}}$,

$$
\mathbf{1}_{V} z_{l} \xrightarrow{\text { dist }} \mathbf{1}_{V} A \otimes G,
$$

which shows that

$$
z_{l} \xrightarrow{\mathrm{wcd}} A \otimes G .
$$

In the non-separable case, use the fact that the $x_{n}$ live in a separable sublattice $X(\Omega, \mathscr{B}, \mu)$ and that $X$ is rearrangement invariant ( $V$ can be put down into a fixed separable superspace $(\Omega, \mathscr{C}, \mu)$ of $(\Omega, \mathscr{A}, \mu)$ by a measure-preserving transformation leaving elements of $\mathscr{B}$ invariant).

Lemma 12. Let $X$ satisfy the hypotheses of Theorem 7. If $l_{2}$ is $C$-representable as a complemented subspace of $X$, then for every $\varepsilon>0$ there is a special $C(1+\varepsilon)$-factorization of the identity of $l_{2}$ through $X$ which maps the basis of $l_{2}$ onto a sequence which is $X$-equiintegrable and converges wcd to a conditionally Gaussian variable.

Proof. (A) Let

$$
l_{2} \xrightarrow{i} X \xrightarrow{\pi} l_{2}
$$

be a factorisation of $\mathrm{id}_{l_{2}}$ through $X$, with $\|\pi\|\|i\| \leq C$. Then $E=i\left(l_{2}\right)$ is a complemented hilbertian subspace of $X$. We show first that we may suppose that the $l_{2}$ basis of $E$ converges wcd to a conditionally Gaussian variable. As at beginning of the proof of Lemma 9, we find in $E$ a normalized sequence $x_{n}=i\left(y_{n}\right)$ which converges wcd to a conditionally gaussian variable. Strictly speaking, this wcd convergence of $\left(x_{n}\right)$ was obtained only on a certain $\mu$-integrable set $U$, such that the norms of $X$ and that of $L_{1}(U)$ are equivalent on $E_{0}$; but, if $\left(\Omega_{p}\right)_{p}$ is an increasing sequence of $\mu$-integrable subsets containing $U$, whose union contains all the supports of the $x_{n}$, we can by a diagonal argument construct the sequence ( $x_{n}$ ) converging wcd on each set $\Omega_{p}$ to a conditionally Gaussian variable. Then the sequence $\left(x_{n}\right)_{n}$ is wcd convergent on the whole of $\Omega$ to a conditionally Gaussian variable. The $y_{n}$ 's may be taken as successive norm one blocks on the basis of $l_{2}$, hence forming
isometrically a $l_{2}$-basis. Let $H=\overline{\operatorname{span}}\left[y_{n}\right]_{n=1}^{\infty}$ and $Q$ be the orthogonal projection of $l_{2}$ onto $H$. Then $\mathrm{id}_{H}=\left.(Q \circ \pi) \circ i\right|_{H}$ gives a $C$-representation of $H$ (itself isometric to $l_{2}$ ) as complemented subspace of $X$, with $l_{2}$ basis $\left(x_{n}\right)$.
(B) Note that every subsequence of $\left(x_{n}\right)$ gives raise to a $C$-representation of $l_{2}$ as complemented subspace of $X$. Thus we may suppose that we have the splitting $x_{n}=x_{n}^{\prime}+x_{n}^{\prime \prime}$ into $X$-equiintegrable and disjoint part. Let $E$ $=\overline{\operatorname{span}}\left[x_{n}\right], E^{\prime}=\overline{\operatorname{span}}\left[x_{n}^{\prime}\right], E^{\prime \prime}=\overline{\operatorname{span}}\left[x_{n}^{\prime \prime}\right]$ and $S^{\prime}: E \rightarrow E^{\prime}$, resp. $S^{\prime \prime}: E \rightarrow E^{\prime \prime}$ the natural operators $\left(S^{\prime} x_{n}=x_{n}^{\prime}, S^{\prime \prime} x_{n}=x_{n}^{\prime \prime}\right.$ ). By the proof of Lemma 9 , we may suppose that ( $x_{n}^{\prime}$ ) converges wcd to a non-zero conditionally Gaussian variable. After extracting a subsequence if necessary, we may suppose (by Lemma 10 and Remark 11) that the unit ball of $E^{\prime}$ is $X$-equiintegrable, and that every weakly null sequence in $E^{\prime}$ converges wcd to a conditionally gaussian variable. As the $x_{n}^{\prime \prime}$ are asymptotically disjoint from the unit ball of $E^{\prime}$, (due to the $X$-equiintegrability of this unit ball), it is easy to see that the restriction of $S^{\prime}$ and $S^{\prime \prime}$ to $E_{n}=\overline{\operatorname{span}}\left[x_{k}\right]_{k \geq n}$ satisfy $\left\|\left.S^{\prime}\right|_{E_{n}}\right\| \leq 1+\varepsilon(n)$, $\left|\left|S^{\prime \prime}\right|_{E_{n}} \| \leq 1+\varepsilon(n)\right.$, with $\varepsilon(n) \rightarrow 0$ as $n \rightarrow \infty$. Given $\varepsilon>0$, we may suppose $\varepsilon(n) \leq \varepsilon$.

Let

$$
l_{2} \xrightarrow{i} X \xrightarrow{\pi} l_{2}
$$

be a $C$-factorization of $\mathrm{id}_{l_{2}}$ through $X$, the image by $i$ of the natural basis $\left(e_{n}\right)$ of $l_{2}$ being $\left(x_{n}\right)$. Let ${ }^{l_{2}} P=i \pi$ be the induced projection from $X$ onto $E=i\left(l_{2}\right)$. Due to Lemma 8, and the hypothesis on $X$, the operator $P S^{\prime \prime}$ is strictly singular. Thus there exists a sequence ( $u_{n}$ ) of successive normalized blocks on the basis $\left(e_{n}\right)$ such that $P S^{\prime \prime} y_{n} \rightarrow 0$ as $n \rightarrow \infty$, where $y_{n}=i\left(u_{n}\right)$. Since $i$ is an isomorphism into $X$, we have in fact $\pi S^{\prime \prime} y_{n} \rightarrow 0$ as $n \rightarrow \infty$. But $u_{n}-\pi S^{\prime} i\left(u_{n}\right)=\pi\left(y_{n}\right)-\pi S^{\prime} y_{n}=\pi S^{\prime \prime} y_{n} \rightarrow 0$, hence (after extraction) we can suppose that the restriction to $H=\overline{\operatorname{span}}\left[u_{n}\right]$ of the operator $I-\pi S^{\prime} i$ is of norm $\leq \varepsilon$. Let $Q$ be the orthogonal projection from $l_{2}$ onto $H$. We have

$$
\left\|\left.\left(I-Q \pi S^{\prime}\right)\right|_{H}\right\|=\left\|\left.\left(Q\left(I-\pi S^{\prime}\right)\right)\right|_{H}\right\| \leq \varepsilon
$$

Then $J=\left.Q \pi S^{\prime} i\right|_{H}$ is invertible, and $\left\|J^{-1}\right\| \leq 1 /(1-\varepsilon)$. Set $i^{\prime}=S^{\prime} J^{-1}$ : $H \rightarrow X$ and $\pi^{\prime}=Q \pi: X \rightarrow H$. We have $\pi^{\prime} i^{\prime}=\mathrm{id}_{H}$ and

$$
\left\|i^{\prime}\right\|\left\|\pi^{\prime}\right\| \leq(1-\varepsilon)^{-1}\left\|S^{\prime}\right\|\|i\|\|\pi\| \leq C(1-\varepsilon)^{-1}(1+\varepsilon)
$$

Finally $\left(z_{n}\right)=\left(S^{\prime} i J^{-1} u_{n}\right)$ provides a $C(1-\varepsilon)^{-1}(1+\varepsilon)$ representation of $l_{2}$ as a complemented subspace of $X$, and by the choice of $\left(x_{n}\right)$, the sequence ( $z_{n}$ ) is $X$-equiintegrable and converges wcd to a conditionally Gaussian variable.

Lemma 13. Under the hypotheses of Theorem 7, if $l_{2}$ is $C$-representable as a complemented subspace of $X$, there exists a sequence $\left(A \otimes G_{n}\right)_{n}$ (of conditionally i.i.d. gaussian variables in $X(\Omega \times[0,1])$ ) whose closed linear span is $C$-complemented in $X$.

Proof. (A) We make first a little digression about complemented spreading models in ultrapowers.

Let $X$ be a Banach space and $\left(x_{k}\right)_{k=1}^{\infty}$ a sequence without converging subsequence. Let $\mathscr{U}$ be a non trivial ultrafilter over the index $\operatorname{set} \mathbf{N}$.

The sequence $\left(x_{k}\right)_{k}$ defines an element $\xi_{1}$ of the ultrapower $\tilde{X}_{1}=X^{\mathrm{N}} / \mathscr{U}$, then an element $\xi_{2}$ of $\tilde{X}_{2}=X_{1}^{\mathrm{N}} / \mathscr{U}, \ldots$ We define recursively a sequence $\left(\tilde{X}_{k}\right)_{k}$ of successive ultrapowers and a sequence $\left(\xi_{k}\right)_{k}, \xi_{k} \in \tilde{X}_{k}$. Thus we have

$$
\begin{aligned}
&\left\|\alpha_{1} \xi_{1}+\cdots+\alpha_{k} \xi_{k}\right\|_{\tilde{X}_{k}}=\lim _{n_{k}, \mathscr{U}} \lim _{n_{k-1}, \mathscr{U}} \cdots \lim _{n_{1}, \mathscr{U}}\left\|\alpha_{1} x_{n_{1}}+\cdots+\alpha_{k} x_{n_{k}}\right\|_{X} \\
& \forall \alpha_{1}, \ldots \alpha_{k} \in \mathbf{R}
\end{aligned}
$$

The whole sequence $\left(\xi_{k}\right)_{k}$ can be considered as living in the same space $\tilde{X}=\Pi_{k} \tilde{X}_{k} / \mathscr{U}$ (which is an ultrapower of $X$ ) and spans in $\tilde{X}$ the so-called spreading model associated to the sequence $\left(x_{k}\right)_{k}$ and the ultrafilter $\mathscr{U}$.

Suppose now that $\left(x_{n}\right)_{n}$ is the image of a 1 -symmetric basis $\left(e_{n}\right)_{n}$ of a reflexive space $Z$ under a $C$-factorization of $\mathrm{id}_{Z}$ through $X$. For each $k$, set $Z_{k}=\operatorname{span}\left[e_{1}, \ldots e_{k}\right]$; for each multiindex $\left(n_{1}, \ldots n_{k}\right) \in \mathbf{N}^{k}$, let $\sigma_{n_{1}, \ldots n_{k}}$ : $Z_{k} \rightarrow Z$ be defined by $\sigma_{n_{1}, \ldots n_{k}}\left(e_{l}\right)=e_{n_{i}}$; set $Z_{n_{1}, \ldots, n_{k}}=\sigma_{n_{1}, \ldots n_{k}}\left(Z_{k}\right)$ and let $r_{n_{1}, \ldots, n_{k}}$ be the natural projection $Z \rightarrow Z_{n_{1}, \ldots, n_{k}}$. Set

$$
\pi_{n_{1}, \ldots n_{k}}=\sigma_{n_{1}, \ldots n_{k}}^{-1} \circ r_{n_{1}, \ldots, n_{k}} \circ \pi \quad \text { and } \quad i_{n_{1}, \ldots n_{k}}=i \circ \sigma_{n_{1}, \ldots n_{k}}
$$

we obtain a factorization $\pi_{n_{1}, \ldots n_{k}}{ }^{\circ} i_{n_{1}, \ldots n_{k}}$ of $\operatorname{id}_{Z_{k}}$ through span $\left[x_{n_{1}}, \ldots x_{n_{k}}\right]$ with $\left\|\pi_{n_{1}, \ldots n_{k}}\right\|\left\|i_{n_{1}, \ldots n_{k}}\right\| \leq C$. Thus we have a $C$-factorization $\tilde{\pi}_{k} \circ \tilde{i}_{k}$ of $\mathrm{id}_{Z_{k}}$ through $\tilde{X}_{k}$ by setting

$$
\tilde{\boldsymbol{i}}_{k} z=\text { class of the family }\left(i_{n_{1}, \ldots n_{k}} z\right)_{\left(n_{1}, \ldots n_{k}\right) \in \mathbf{N}^{k}}
$$

and

$$
\tilde{\pi}_{k} \tilde{x}=\lim _{n_{k}, \mathscr{U} n_{k-1}, \mathscr{U}} \lim _{n_{1}, \mathscr{U}} \ldots \lim _{n_{1}} \pi_{n_{1}, \ldots, n_{k}}\left(x_{n_{1}, \ldots n_{k}}\right)
$$

when $\tilde{x}$ is the class of the family $\left(x_{n_{1}, \ldots n_{k}}\right)_{\left(n_{1}, \ldots n_{k}\right) \in \mathbf{N}^{k}}$. (Here the limits are norm limits in finite dimensional spaces.) The image by $\tilde{\pi}_{k}$ of the basis of $Z_{k}$ is the sequence $\left(\xi_{1}, \ldots \xi_{k}\right)$. We thus obtain a $C$-factorization $\tilde{\pi} \circ \tilde{\imath}$ of $\mathrm{id}_{Z}$ through $\tilde{X}$ by letting $\tilde{\imath} \tilde{z}$ be class of the family $\left(\tilde{\imath}_{k} r_{1,2, \ldots k} z\right)_{k}$ and $\tilde{\pi} \tilde{x}=$
$\lim _{k, \mathscr{U}} \tilde{\pi}_{k} \tilde{x}_{k}$ when $\tilde{x}$ is the class of the family $\left(\tilde{x}_{k}\right)_{k}$, where the limits are taken, say, coordinatewise. The image of the basis of $Z$ by $\tilde{\pi}$ is the "fundamental sequence of the spreading model" $\left(\xi_{n}\right)_{n}$.
(B) Now we use this ultrapower construction starting with the hilbertian $X$-equiintegrable sequence given by Lemma 12 . In this case the elements $\left(\xi_{k}\right)_{k=1}^{\infty}$ (constructed above) live in the band $\tilde{X}_{e q}$ of the ultrapower $\tilde{X}$ whose elements are defined by $X$-equiintegrable families of elements of $X$. As is well known, this band is nothing but a space $X(\tilde{\Omega}, \tilde{\mathscr{A}}, \tilde{\mu})$, where $(\tilde{\Omega}, \tilde{\mathscr{A}}, \tilde{\mu})$ is (a Stone representation of) the ultrapower of the measure space $(\Omega, \mathscr{A}, \mu)$ (or, equivalently, $\left.L_{1}(\tilde{\Omega}, \tilde{\mathscr{A}}, \tilde{\mu})=\left[L_{1}(\Omega, \mathscr{A}, \mu)\right]_{e q}^{\sim}\right)$. Moreover, the distribution of the sequence $\left(\xi_{k}\right)_{k}$ is the limit of that of the $\left(x_{k}\right)$; to be more explicit,

$$
\left(x, x_{n_{1}}, x_{n_{2}}, \ldots x_{n_{k}}\right) \frac{\text { dist }}{n_{1}, \mathscr{U} ; n_{2}, \mathscr{U} ; \ldots n_{k}, \mathscr{U}}\left(x, \xi_{1}, \ldots \xi_{k}\right) \quad \forall x \in X
$$

where the convergence in distribution is evaluated against bounded continuous functions on $\mathbf{R}^{k+1}$ [DC]. But we have

$$
\left(x_{n_{1}}, x_{n_{2}}, \ldots x_{n_{k}}\right) \xrightarrow[n_{1} \rightarrow \infty ; n_{2} \rightarrow \infty ; \ldots n_{k} \rightarrow \infty]{\mathrm{wcd}}\left(Y_{1}, Y_{2}, \ldots Y_{k}\right)
$$

where the $Y_{j} \in L_{1}(\Omega \times S)$ are conditionally Gaussian, independent and identically distributed. Hence the sequence $\left(\xi_{k}\right)_{k}$ is conditionally equivalent in distribution to a sequence of conditionally independent equidistributed Gaussian variables, i.e.,

$$
\forall x \in X, \forall k \in \mathbf{N}:\left(x, \xi_{1}, \ldots \xi_{k}\right) \stackrel{\text { dist }}{\sim}\left(x, A \otimes G_{1}, A \otimes G_{2}, \ldots A \otimes G_{k}\right)
$$

where $A \in L_{0}(\Omega, \mathscr{A}, \mu)$ and $\left(G_{k}\right)_{k}$ is a sequence of independent normal Gaussian variables, defined on ( $[0,1], \lambda$ ). We can suppose that $A>0$ on a subset $U$ of $\Omega$. Let $\mathscr{B}_{k}$ be the $\lambda$ complete $\sigma$-algebra generated by the variable $G_{k}$ and $\mathscr{C}$ that generated by $\mathscr{A}_{U}$ (the trace of $\mathscr{A}$ on $U$ ) and the variables $\xi_{1}, \ldots \xi_{k}$. There is an isomorphism $\mathscr{T}$ of measure algebras from ( $\mathfrak{b}, \tilde{\mu}$ ) onto

$$
\left(\mathscr{A}_{U} \bar{\otimes} \mathscr{B}_{1} \ldots \bar{\otimes} \mathscr{B}_{k} \bar{\otimes} \ldots, \mu \otimes \lambda \ldots \otimes \lambda \ldots\right)
$$

which is the identity on $\mathscr{A}$ and maps $\xi_{j}$ on $A \otimes G_{j}$. This isomorphism $\mathscr{T}$ generates a lattice isometry from the space $X(\mathscr{b}, \mu)$ onto

$$
X\left(U \times[0,1]^{\mathbf{N}}, \mathscr{A}_{U} \otimes \bigotimes_{k} \mathscr{B}_{k}, \mu \otimes \lambda^{\otimes \mathbf{N}}\right)
$$

hence the sequence $\left(Y_{k}\right)_{k}=\left(A \otimes G_{k}\right)_{k}$ spans a complemented closed space
in the last space. By a standard isomorphism theorem, we may replace $[0,1]^{\mathrm{N}}$ by $[0,1]$.
(C) More precisely we have $Y_{k}=\pi\left(e_{k}\right)$ for a certain $C(1+\varepsilon)$-factorization $\pi \circ i$ of $\operatorname{id}_{l_{2}}$ through $X(\Omega \times[0,1])$. As $\left(Y_{n}\right)$ is isometrically equivalent to the basis of $l_{2}$, this means that $\operatorname{span}\left[Y_{n}\right]_{n}$ is $C(1+\varepsilon)$-complemented in $X(\Omega \times$ $[0,1])$. A more careful treatment would show that the tail $\left(z_{n}\right)_{n \geq m}$ of the sequence $\left(z_{n}\right)_{n}$ constructed in the proof of Lemma 12 provides in fact a $C\left(1+\varepsilon_{m}\right)$-representation of $l_{2}$ as complemented subspace of $X$, with $\varepsilon_{m} \rightarrow 0$ as $m \rightarrow \infty$; as a consequence, the sequence $\left(Y_{l}\right)_{l}$ spans in fact a $C$-complemented space.

Proof of Theorem 7. Theorem 7 is now an immediate consequence of Lemma 13 and Proposition 6.

## 4. On the isomorphism $X \simeq X\left(l_{2}\right)$ for Orlicz spaces

It is well known (see [LT2], Prop. 2d4), that a separable r.i. function space $X$ with nontrivial Boyd indices is isomorphic (as a Banach space) to its vectorial extension $X\left(l_{2}\right)$. In this section we show that the converse statement is true for Orlicz function spaces.

First we fix some notations. If $X$ is a Banach lattice of measurable functions on $(\Omega, \mathscr{A}, \mu)$, then $X\left(l_{2}\right)$ denotes the space of $\mathscr{A}$-measurable functions $f: \Omega \rightarrow l_{2}\left(l_{2}\right.$ being equipped with its Borel $\sigma$-field) for which the scalar function $\omega \stackrel{\mapsto}{\mapsto}\|(\omega)\|_{l_{2}}$ belongs to $X$, while $X\left[l_{2}\right]$ is the closure in $X\left(l_{2}\right)$ of the algebraic tensor product $X \otimes l_{2}$ (consisting of finite sums $\left.\sum_{i=1}^{n} f_{i} \otimes x_{i}, f_{i} \in X, x_{i} \in l_{2}\right)$. Our $X\left[l_{2}\right.$ ] is denoted by $X\left(l_{2}\right)$ in [LT2], while, when $X=X^{\prime \prime}$, the space $X\left(l_{2}\right)$ coincides with that denoted by $\overline{X\left(l_{2}\right)}$ in [LT2]. When $X$ is order continuous, $X\left(l_{2}\right)=X\left[l_{2}\right]$.

If $\varphi: \mathbf{R}_{+} \rightarrow \mathbf{R}_{+}$is an Orlicz function, i.e., a convex increasing function with $\varphi(0)=0$, then $L_{\varphi}(\Omega, \mathscr{A}, \mu)$ is, as usual, the space of measurable functions $f$ : $\Omega \rightarrow \mathbf{R}$ such that

$$
\varphi(c|f|) \in L_{1}(\Omega, \mathscr{A}, \mu) \quad \text { for some } c>0
$$

and we denote by $M_{\varphi}$ the closure in $L_{\varphi}$ of the space of simple $\mu$-integrable functions.

We have the following result.
Theorem 14. The following assertions are equivalent:
(i) $L_{\varphi}$ is reflexive.
(ii) $L_{\varphi}$ is isomorphic to $L_{\varphi}\left[l_{2}\right]$.
(iii) $L_{\varphi}$ is isomorphic to $L_{\varphi}\left(l_{2}\right)$.
(iv) $M_{\varphi}$ is isomorphic to $M_{\varphi}\left(l_{2}\right)$.

This is a consequence of the following result.

Proposition 15. An Orlicz function space never contains a complemented subspace isomorphic to $l_{1}\left(l_{2}\right)$.

Proof. Suppose that $X=L_{\varphi}$ contains $l_{1}\left(l_{2}\right)$ as complemented subspace. Then its bidual $X^{* *}$ contains $\left(l_{1}\left(l_{2}\right)\right)^{* *}=\left(l_{\infty}\left(l_{2}\right)\right)^{*}$ as complemented subspace. The structure of the dual of an Orlicz space is well known (see [An] for the case where the measure space is finite, [Fer] for the general case). We have:

$$
L_{\varphi}(\Omega, \mathscr{A}, \mu)^{*}=L_{\varphi_{*}}(\Omega, \mathscr{A}, \mu) \oplus L_{1}(S, \Sigma, \sigma)
$$

where $\varphi_{*}$ is the Young conjugate of $\varphi$. (The $L_{1}$-component is null if $\varphi$ verifies the $\Delta_{2}$-condition; in the other case it is non-separable). Thus the bidual is given by

$$
X^{* *}=L_{\varphi}(\Omega, \mathscr{A}, \mu) \oplus L_{1}(T, \mathscr{T}, \tau) \oplus L_{\infty}(S, \Sigma, \sigma)
$$

where again $L_{1}(T)$ is either null (if $\varphi$ satisfies the $\Delta_{2}^{*}$-condition) or non-separable. On the other hand $\left(l_{\infty}\left(l_{2}\right)\right)^{*}$ contains $l_{\infty}^{*}\left(l_{2}\right)$ as complemented subspace. Here $l_{\infty}^{*}\left(l_{2}\right)$ can be defined in an abstract way by using Krivine functional calculus (see [LT2], 1d1), but in fact it identifies with the Banach space $M$ whose elements are the sequences $\left(\mu_{n}\right)$ in $\left(l_{\infty}\right)^{*}$ such that

$$
\left\|\left(\mu_{n}\right)\right\|:=\operatorname{Sup}\left\{\sum_{i=1}^{N}\left\langle\mu_{n}, x_{n}\right\rangle / N \geq 1, x_{n} \in l_{\infty}, \forall n=1, \ldots N ;\left\|\sum_{n=1}^{N} x_{n}^{2}\right\|_{\infty} \leq 1\right\}
$$

is finite. If $\mu_{n}=0$ for $n>N$, then $\left(\mu_{n}\right)_{n}$ identifies to $\sum_{i=1}^{N} \mu_{n} \otimes e_{n}$ ( $\left(e_{n}\right)$ is the natural basis of $\left.l_{2}\right)$. This space in turn isometrically identifies to a subspace of $\left(l_{\infty}\left(l_{2}\right)\right)^{*}$ by setting

$$
\left\langle\left(\mu_{n}\right)_{n},\left(x_{i, n}\right)_{i, n}\right\rangle=\sum_{n}\left\langle\mu_{n},\left(x_{i, n}\right)_{i}\right\rangle
$$

for every $\left(\mu_{n}\right)_{n} \in M$ and $\left(x_{i, n}\right)_{i, n} \in l_{\infty}\left(l_{2}\right)$. If $F \in\left(l_{\infty}\left(l_{2}\right)\right)^{*}$, we define $P F \in$ $l_{\infty}^{*}\left(l_{2}\right)$ by $P F=\left(\mu_{n}\right)_{n}$ where $\forall x \in l_{\infty},\left\langle\mu_{n}, x\right\rangle=\left\langle F, x \otimes e_{n}\right\rangle$; then $P$ is a natural projection from $\left(l_{\infty}\left(l_{2}\right)\right)^{*}$ onto $l_{\infty}^{*}\left(l_{2}\right)$.

Note that $l_{\infty}^{*}$ is a $L_{1}$-space containing a 1 -complemented sublattice isomorphic to a $l_{1}(\Gamma)$ space, defined on an index set $\Gamma$ which has cardinality $2^{c}$, where $\mathfrak{c}$ is the cardinality of the continuum $(\Gamma=\beta \mathbf{N})$. We obtain thus that $l_{1}(\Gamma)\left(l_{2}\right)$ embeds isomorphically as a complemented subspace $E$ in $X^{* *}=$ $L_{\varphi} \oplus L_{1}(T) \oplus L_{\infty}(S)$.

Let $\left(e_{\gamma, n}\right)_{\gamma \in \Gamma, n \in \mathbf{N}}$ be the $l_{1}(\Gamma)\left(l_{2}\right)$ basis in $X^{* *}$. For each $\gamma \in \Gamma$, let $E_{\gamma}=\overline{\operatorname{span}}\left[e_{\gamma, n}\right]_{n \geq 1}$. We also denote by $Q$ the given projection $X^{* *} \rightarrow E$, and by $P_{1}$ (resp. $P_{\infty}, P_{0}$ ) the natural projections $X^{* *} \rightarrow L_{1}(T)$ (resp. $X^{* *} \rightarrow$ $\left.L_{\infty}(S), X^{* *} \rightarrow L_{\varphi}(\Omega)\right)$. For all $\varepsilon>0$, there exists $y_{\gamma} \in E_{\gamma}$ with $\left\|y_{\gamma}\right\|=1$ and
$\left\|Q P_{1} y_{\gamma}\right\|<\varepsilon,\left\|Q P_{\infty} y_{\gamma}\right\|<\varepsilon$ : if not, e.g., $Q P_{1}$ would be an isomorphism into. In particular $P_{1}: E_{\gamma} \rightarrow P_{1} E_{\gamma}$ and $Q: P_{1}\left(E_{\gamma}\right) \rightarrow Q P_{1}\left(E_{\gamma}\right)$ are isomorphisms; let $J$ be the inverse isomorphism of the last one.
$Q P_{1}\left(E_{\gamma}\right)$ is a hilbertian subspace of $E$, where $E$ is isomorphic to $l_{1}\left(l_{2}\right)$ : hence it contains a further space $Z$, which is complemented in $E$, by a projection $\pi$. Then $J Z$ would be a hilbertian subspace of $P_{1}\left(E_{\gamma}\right)$, hence of $L_{1}(T)$, and would be complemented in $X^{* *}$, and a fortiori in $L_{1}(T)$, by the projection $J \pi Q$; this is impossible. The same reasoning works with $P_{\infty}$ in place of $P_{1}$.

Then we have

$$
\forall\left(\alpha_{\gamma}\right)_{\gamma} \in \mathbf{R}^{(\Gamma)},\left\|Q\left(P_{1}+P_{\infty}\right)\left(\sum_{\gamma} \alpha_{\gamma} y_{\gamma}\right)\right\| \leq \varepsilon \sum_{\gamma}\left|\alpha_{\gamma}\right| \leq \varepsilon C\left\|\sum_{\gamma} \alpha_{\gamma} y_{\gamma}\right\|
$$

where $C$ is the equivalence constant of $E$ with $l_{1}\left(l_{2}\right)$. Hence

$$
\left\|Q P_{0}\left(\sum_{\gamma} \alpha_{\gamma} y_{\gamma}\right)\right\| \geq(1-\varepsilon C)\left\|\sum_{\gamma} \alpha_{\gamma} y_{\gamma}\right\|
$$

(since $\left.Q y_{\gamma}=y_{\gamma}, \forall \gamma\right)$ ). Then

$$
\left\|\sum_{\gamma} \alpha_{\gamma} P_{0} y_{\gamma}\right\|=\left\|P_{0}\left(\sum_{\gamma} \alpha_{\gamma} y_{\gamma}\right)\right\| \geq(1-\varepsilon C)\|Q\|^{-1}\left\|\sum_{\gamma} \alpha_{\gamma} y_{\gamma}\right\| ;
$$

i.e., $\left(P_{0} y_{\gamma}\right)_{\gamma}$ spans a subspace of $L_{\varphi}(\Omega)$ isomorphic to $l_{1}(\Gamma)$. If $(\mathscr{A}, \mu)$ is countably generated as measure algebra, this is impossible, since the density character of $L_{\varphi}(\Omega)$ is at most $c$, while that of $l_{1}(\Gamma)$ is $2^{c}$. In the general case, we can find a sub- $\sigma$ algebra $\mathscr{B}$ of $\mathscr{A}$, such that $(\mathscr{B}, \mu)$ is countably generated, and such that the elements of the $\left(l_{1}\left(l_{2}\right)\right)$-basis in $L_{\varphi}(\Omega, \mathscr{A}, \mu)$ are $\mathscr{B}$-measurable.

Remark. The same result is true (with same proof) for $M_{\varphi}$ in place of $L_{\varphi}$.
Proof of Theorem 14. If $L_{\varphi}$ is reflexive, then $L_{\varphi}=M_{\varphi}$ and $L_{\varphi}\left(l_{2}\right)=$ $L_{\varphi}\left(l_{2}\right)=M_{\varphi}\left(l_{2}\right)$, hence the assertions (ii)-(iv) are the same, and in fact a consequence of [LT2], prop. 2 d 4.

Conversely, suppose that one of the conditions (ii), (iii), or (iv) is verified. We show first that $\varphi$ verifies the condition $\Delta_{2}^{*}$ (i.e., is equivalent to a $p$-convex Orlicz function, for some $p>1$ ). If not, then by [L] there exists a complemented subspace $F$ of $L_{\varphi}$, (resp. $M_{\varphi}$ ) isomorphic to $l_{1}$ and spanned by disjoint positive functions $\left(f_{i}\right)_{i=1}^{\infty}$ (in fact, indicator functions). There is a positive projection $P$ with range $F$. If $\left(e_{i}\right)_{i}$ denotes the natural basis of $l_{2}$, then $\left(f_{i} \otimes e_{i}\right)_{i}$ span in $L_{\varphi}\left[l_{2}\right]$ (resp. $\left.M_{\varphi}\left(l_{2}\right)\right)$ a subspace $E$ isomorphic to
$l_{1}\left(l_{2}\right)$, and $P \otimes \mathrm{id}_{l_{2}}$ defines a projection from $L_{\varphi}\left(l_{2}\right)$ (resp $\left.M_{\varphi}\left(l_{2}\right)\right)$ onto $E$. The boundedness of $P \otimes \mathrm{id}_{l_{2}}$ is a consequence of the positivity of $P$ : if $\left(h_{i}\right)_{i} \subset L_{0}$ with $\left(\sum_{i}\left|h_{i}\right|^{2}\right)^{1 / 2} \in L_{\varphi}$, then
hence $l_{1}\left(l_{2}\right)$ appears as a complemented subspace of $L_{\varphi}\left(\operatorname{resp} M_{\varphi}\right)$, which is impossible.

If now $\varphi$ does not verify the $\Delta_{2}$-condition, then $M_{\varphi}$ contains $c_{0}$ as complemented sublattice spanned by disjoint elements (again indicator functions) and $L_{\varphi}$ contains $l_{\infty}$ as complemented sublattice; again with positive projection. Then:

$$
c_{0}\left(l_{2}\right) \subset_{c} M_{\varphi}\left(l_{2}\right), \quad l_{\infty}\left[l_{2}\right] \subset_{c} L_{\varphi}\left[l_{2}\right], \quad l_{\infty}\left(l_{2}\right) \subset_{c} L_{\varphi}\left(l_{2}\right)
$$

and we deduce

$$
c_{0}\left(l_{2}\right){\underset{c}{ }}^{M_{\varphi}}, \quad \text { or } \quad l_{\infty}\left[l_{2}\right] \subset_{c} L_{\varphi}, \quad \text { or } \quad l_{\infty}\left(l_{2}\right) \subset_{c} L_{\varphi}
$$

Dualizing and using the fact that $l_{\infty}^{*}\left(l_{2}\right) \subset_{c}\left(l_{\infty}\left[l_{2}\right]\right)^{*}$ and $\subset_{c}\left(l_{\infty}\left(l_{2}\right)\right)^{*}$, we obtain

$$
l_{1}\left(l_{2}\right) \underset{c}{\subset} M_{\varphi}^{*}=L_{\varphi_{*}} \quad \text { or } \quad l_{\infty}^{*}\left(l_{2}\right) \subset c_{c} L_{\varphi}^{*}=L_{\varphi_{*}} \oplus L_{1}(S)
$$

The first assertion is impossible by Prop. 15, and the second one also, by the proof of Prop. 15.
5. On the isomorphism $X\left(l_{2}\right) \approx X$ for $q$-concave r.i. spaces $(q<2)$

The main result of this section is the following:
Theorem 16. Let $X$ be a $q$-concave $(q<2)$ rearrangement invariant space over $\Omega=[0,1]$. Then $X\left(l_{2}\right)$ is isomorphic to $X$ iff $X$ has non-trivial lower Boyd index $\left(p_{X}>1\right)$.

The following criterion will be used:
Lemma 17. Let $Y$ be an r.i. space (over $\Omega=[0,1]$ or $[0, \infty]$ ); $G$ a normalized Gaussian variable (defined on the probability space ( $S, \Sigma, \mathbf{P}$ )). Denote by $Z$ the space of measurable functions $f \in L_{0}(\Omega)$ such that $f \otimes G$
belongs to $Y^{\prime \prime}(\Omega \times S)$. Then $Y$ has nontrivial upper Boyd index iff it is (algebraically) included in $Z$.

Proof. (a) Suppose that $Y$ has non trivial upper Boyd index $q_{Y}<\infty$. Let

$$
U_{k}=\{k \leq|G|<k+1\}
$$

Then $\mathbf{P}\left(U_{k}\right) \leq P[|G| \geq k] \leq C \cdot e^{-k^{2} / 2}$. Let $f \in Y$; then $\forall \rho \leq 1, \forall q>q_{Y}$, $\left\|D_{\rho} f\right\|_{Y} \leq C_{q} \rho^{1 / q}\|f\|_{Y}$ (where $D_{\rho}$ is the usual dilation operator: $D_{\rho} f(t)=$ $f(t / \rho)$ ). Hence

$$
\left\|f \otimes \mathbf{1}_{U_{k}}\right\|_{Y(\Omega \times S)} \leq C e^{-k^{2} / 2 q}\|f\|_{Y(\Omega)}
$$

and $\sum_{k} f \otimes\left(\mathbf{1}_{U_{k}} G\right)$ converges in $Y(\Omega \times S)$; i.e., $f \otimes G \in Y(\Omega \times S)$.
Conversely suppose that $\|f\|_{Z} \leq C\|f\|_{Y}$ for all $f \in Y$. To prove that $q_{Y}<\infty$, it suffices to prove that $\left\|D_{a}\right\|_{Y \rightarrow Y}<1$ for some $a>0$. But setting $a_{k}=\mathbf{P}\left(U_{k}\right)$,

$$
\left\|D_{a_{k}} f\right\|_{Y}=\left\|f \otimes \mathbf{1}_{U_{k}}\right\|_{Y} \leq \frac{1}{k}\|f \otimes G\|_{Y} \leq \frac{C}{k}\|f\|_{Y}
$$

hence for sufficiently large $k$, we are done.
To prove Theorem 16, we are led to use the following proposition, the proof of which is inspired from that of Theorem 5.6 of [JMST], relative to r.i. spaces (over [0,1]) which embed complementably in a given $q$-concave r.i. function space. Recall that a quasi-norm $\|\|$ on a space $X$ verifies $\left\|x_{1}+x_{2}\right\| \leq \gamma\left(\left\|x_{1}\right\|+\left\|x_{2}\right\|\right)$ for some $\gamma>0$ and all $x_{1}, x_{2} \in X$.

Proposition 18. Let $X$ be a $\frac{1}{2}$-convex, 1-concave quasi-Banach r.i. function space, over [0, 1], Y a (Banach) r.i. function space (over [0, 1]). We suppose that the 2-convexified spaces $X^{(2)}$ and $Y^{(2)}$ are in duality $\left(X^{(2)}\right.$ is the Köthe dual of $Y^{(2)}$ ). Let $Z$ be another (Banach) r.i. function space, algebraically included in Y. Suppose that there exist positively linear bounded operators

$$
T: Y_{+} \rightarrow Z_{+} \quad \text { and } \quad S: X_{+} \rightarrow X_{+}
$$

such that for every $f$ in $Y_{+}$and $g$ in $X_{+},\left\langle(T f)^{1 / 2},(S g)^{1 / 2}\right\rangle \geq\left\langle f^{1 / 2}, g^{1 / 2}\right\rangle$, and $T$ is order continuous. Then $Y=L_{1}$ or $Y=Z$.

Before proving Prop. 18, we state several lemmas.
Lemma 19. Let $X(\Omega)$ be a concave (quasi-Banach) Köthe function space over a finite measure space $(\Omega, \mathscr{A}, \mu)$, containing algebraically $L_{1}(\Omega)$, and $T$ : $L_{1}^{+}(S) \rightarrow X_{+}(\Omega)$ be a positively linear bounded operator. Then there exists a
measurable $\mu$-a.e. positive function $\psi$ defined on $\Omega$ such that $T$ factorizes through $L_{1}(\psi)=L_{1}(\Omega, \mathscr{A}, \psi \cdot \mu)$, i.e.,

$$
\forall f \in L_{1}^{+}(S), \quad \int \psi \cdot T f d \mu \leq C \int f d \mu
$$

and

$$
\forall g \in L_{1}^{+}(\psi), \quad\|g\|_{X} \leq C \int \psi \cdot g d \mu
$$

We thus have $T=i \circ \tilde{T}$, where $i: L_{1}(\psi) \rightarrow X$ is the identity map and $\tilde{T}$ : $L_{1}^{+} \rightarrow L_{1}(\psi)$ acts as $T$.

Lemma 20. Let $Y=Y(\Omega)$ be a (Banach) Köthe function space over $(\Omega, \mathscr{A}, \mu)$, and $T: Y_{+}(\Omega) \rightarrow L_{1}^{+}(S)$ be an order continuous positively linear bounded operator. There exists an element $\varphi$ of $Y_{+}^{\prime}$ such that $T$ factorizes through $L_{1}^{+}(\varphi)$, i.e.,

$$
\forall f \in Y_{+}, \quad\|T f\|_{1} \leq C \cdot \int \phi \cdot f d \mu \leq C^{\prime}\|f\|_{Y}
$$

We have thus $T=\tilde{T} \circ j$, where $j: Y \rightarrow L_{1}(\varphi)$ is the identity map, and $\tilde{T}$ : $L_{1}^{+}(\Omega) \rightarrow L_{1}^{+}$is formally the same operator as $T$ (i.e., coincides with $T$ on $\left.Y_{+}(\Omega)\right)$.

Lemma 19 and 20 are special cases of Krivine factorization theorem, and the proof of Lemma 19 is close to that of Lemma 5.7 of [JMST].

To obtain Lemma 19 , we suppose $\|T\| \leq 1-\varepsilon$ and separate in the space $L_{\infty}(\Omega)$ the convex sets $L_{\infty}^{+}(\Omega)$ (which has non empty interior) and $C_{1}-C_{2}$, where
$C_{1}=\left\{f \in L_{\infty}^{+}(\Omega) / f \leq T z+h, z \in L_{1}^{+}(S),\|z\|_{1} \leq 1, h \in L_{1}^{+}(\Omega),\|h\|_{1} \leq \varepsilon\right\}$
and

$$
C_{2}=\left\{f \in L_{\infty}^{+}(\Omega) \cap X_{+} /\|f\|_{X}>\gamma\right\}
$$

( $\gamma$ is the quasi-norm constant), by a positive element of $L_{\infty}^{*}$. To obtain Lemma 20, the same reasoning works, if $\|T\| \leq 1$, now with $C_{1}=L_{\infty}^{+}(\Omega) \cap$ $B_{Y}\left(B_{Y}\right.$ is the unit ball of $\left.Y\right)$ and $C_{2}=\left\{f \in L_{\infty}^{+}(\Omega) \cap Y_{+} /\|T f\|_{1}>1\right\}$. We use the order continuity of $T$ to pass from the case $f \in L_{\infty}^{+}(\Omega) \cap Y_{+}$to the case $f \in Y_{+}$.

Lemma 21. Let $Y, Z$ be r.i. (Banach) function spaces over [0,1], with $Z \neq L_{\infty}([0,1])$. Suppose that there exists a positively linear bounded operator $T$ : $Y_{+} \rightarrow Z_{+}$. For all $n \geq 1$ and $i, 1 \leq i \leq 2^{n}$, set $x_{n, i}=\mathbf{1}_{\left[(i-1) \cdot 2^{-n}, i \cdot 2^{-n}\right]}$ and $y_{n}=\max _{1 \leq i \leq 2^{n}} T x_{n, i}$. Suppose $\inf _{n}\left\|y_{n} 1_{\left\{y_{n} \leq R\right\}}\right\|_{Z}>0$ for some $R$. Then the $Y$-norm dominates the $Z$-norm; i.e., there exists a constant $V$ such that $\forall f \in Y \cap Z,\|f\|_{Z} \leq V\|f\|_{Y}$.

This lemma is the positive analog of Lemma 5.2 of [JMST], with plainly analogous proof.

Finally we shall use the following positive version of Thm. 2.1 of [JMST].
Lemma 22. Let $Z$ be an $m$-concave Banach lattice (with constant M). For every finite positively $K$-symmetric sequence $\left(y_{i}\right)_{i=1}^{n}$ in $Z_{+}$, and every choice of positive scalars $\left(a_{i}\right)_{i=1}^{n}$, we have

$$
\begin{aligned}
\frac{1}{D}\left\|\sum_{i=1}^{n} a_{i} y_{i}\right\| & \leq \max \left[\left(\mathbf{E}_{\pi}\left\|\max _{i=1}^{n}\left|a_{\pi(i)} y_{i}\right|\right\|^{m}\right)^{1 / m}, \frac{1}{n}\left\|\sum_{i=1}^{n} y_{i}\right\|_{Z}\left(\sum_{i=1}^{n} a_{i}\right)\right] \\
& \leq K\left\|\sum_{i=1}^{n} a_{i} y_{i}\right\|
\end{aligned}
$$

(where $D=D(K, m, M)$ does not depend on the sequence $\left.\left(y_{i}\right)_{i}\right)$.
(By $\mathbf{E}_{\pi}$ we mean the average for $\pi$ belonging to the group $\mathbf{S}_{n}$ of permutations over $\{1 \ldots, n\}$; the sequence $\left(y_{i}\right)_{i}$ is said to be positively $K$-symmetric if there exists $C$ such that for every nonnegative reals $a_{i}, i=1, \ldots n$, and each permutation $\pi \in \mathbf{S}_{n}$, we have $\left.\left\|\sum a_{i} y_{\pi(i)}\right\| \leq K\left\|\sum a_{i} y_{i}\right\|\right)$.

This lemma can be proven following the method of the proof of Thm. 2.1 of [JMST], or it can also be formally deduced by applying this theorem in $L_{2}\left(Z^{(2)}\right)$, where $Z^{(2)}$ is the 2-convexification of $Z$, to the sequence ( $\varepsilon_{i} \otimes y_{i}^{1 / 2}$ ), where the $\varepsilon_{i}$ are independent Bernouilli variables, and using MaureyKhintchine inequalities [LT2, 1.d.6].

Proof of Proposition 18. (A) The first step is an interpolation procedure, as in the proof of Thm. 5.6 in [JMST] (but we cannot dualize the operators now).

Let $j_{X}$ be the natural injection from $L_{1}([0,1])$ onto $X$. We apply Lemma 19 to the operator $S \circ j_{X}: L_{1}^{+} \rightarrow X$ and obtain a measurable function $\psi>0$ (a.e. on $[0,1])$ such that $S: L_{1}^{+} \rightarrow L_{1}^{+}(\psi)$ is bounded and the identity map: $j_{\psi, X}: L_{1}^{+}(\psi) \rightarrow X$ is also bounded. This implies that the identity map: $L_{2}(\psi) \rightarrow X^{(2)}$ is bounded, and so is its conjugate $Y^{(2)} \rightarrow L_{2}(1 / \psi)$; hence the identity map $i_{Y, 1 / \psi}: Y \rightarrow L_{1}(1 / \psi)$ is bounded. A fortiori the identity map $i_{Z, 1 / \psi}: Z \rightarrow L_{1}(1 / \psi)$ is bounded. We apply now Lemma 20 to the operator
$i_{Z, 1 / \psi} \circ T: Y_{+} \rightarrow L_{1}^{+}(1 / \psi)$ and obtain a measurable function $\varphi \geq 0$ on $[0,1]$ such that $T: L_{1}^{+}(\varphi) \rightarrow L_{1}^{+}(1 / \psi)$ is bounded, and the same for the identity map: $i_{Y, \varphi}: Y \rightarrow L_{1}(\varphi)$. There exists a measurable subset $E$ of $[0,1]$, such that $\varphi$ is bounded from above and below (i.e., $1 / M \leq \varphi \leq M$ ) on $E$. The spaces $Y(E)$ and $X(E)$ (consisting of functions of $Y$, resp. $X$, with support in $E$ ) are lattice isomorphic to $Y$, resp. $X$, by the same dilation operator $D$ (if $|E|=a$ we may suppose $E=[0, a]$ and take $D f(t)=f(a t), \forall t \in[0,1])$. That $D$ is bounded on $X(E)$ is a consequence of the $\frac{1}{2}$-convexity of $X$. This same operator takes $L_{1}(\varphi)(E) \approx L_{1}(E)$ onto $L_{1}([0,1])$. Replacing $T$ and $S$ by $a^{-1 / 2} T \circ D^{-1}$, resp. $a^{-1 / 2} S \circ D^{-1}$ (this normalization conserves the "duality inequality"), we may suppose that $E=[0,1]$ and $L_{1}(\varphi)=L_{1}([0,1])$. We have then the following commutative diagrams:


Fix $0<\theta<1$; we introduce the Calderón-Lozanovski interpolation spaces:

$$
\begin{array}{lll}
Y_{\theta}=\left[Y, L_{1}\right]_{\theta}, & Z_{\theta}=\left[Z, L_{1}(1 / \psi)\right]_{\theta}, & \bar{Y}_{\theta}=\left[Y, L_{1}(1 / \psi)\right]_{\theta} \\
X_{\theta}=\left[X, L_{1}\right]_{\theta}, & \bar{X}_{\theta}=\left[X, L_{1}(\psi)\right]_{\theta} . &
\end{array}
$$

As is well known (this is an easy application of the formula $x^{\theta} y^{1-\theta}=$ $\left.\inf _{t>0}\left((1-\theta) t^{\theta} x+\theta t^{-(1-\theta)} y\right)\right) T$ is bounded as an operator from $Y_{\theta}^{+}$into $Z_{\theta}^{+}$, and so is $S: X_{\theta}^{+} \rightarrow \bar{X}_{\theta}^{+}$.

The 2-convexified spaces $X_{\theta}^{(2)}$ and $Y_{\theta}^{(2)}$ are in duality; for, we have $X_{\theta}^{(2)}=\left[X, L_{1}\right]_{\theta}^{(2)}=\left[X^{(2)}, L_{2}\right]_{\theta}$ and similarly $Y_{\theta}^{(2)}=\left[Y^{(2)}, L_{2}\right]_{\theta}$, and by [Lo],

$$
Y_{\theta}^{(2) \prime}=\left[Y^{(2)}, L_{2}\right]_{\theta}^{\prime}=\left[Y^{(2) \prime}, L_{2}^{\prime}\right]_{\theta}=\left[X^{(2)}, L_{2}\right]_{\theta}=X_{\theta}^{(2)}
$$

Similarly $\bar{X}_{\theta}^{(2)}$ and $\bar{Y}_{\theta}^{(2)}$ are in duality. The inequality $\left\langle(T f)^{1 / 2},(S g)^{1 / 2}\right\rangle \geq$ $\left\langle f^{1 / 2}, g^{1 / 2}\right\rangle$ extends to the case where $f \in Y_{\theta}^{+}, g \in X_{\theta}^{+}$: use the fact that

$$
T f=T\left(\sup \left\{f^{\prime} \mid f^{\prime} \text { simple, } 0 \leq f^{\prime} \leq f\right\}\right) \geq \sup \left\{T f^{\prime} \mid f^{\prime} \text { simple, } 0 \leq f^{\prime} \leq f\right\}
$$

and similarly for $S g$.

We deduce that $\|T f\|_{z_{\theta}} \geq \delta_{\theta}\|f\|_{Y_{\theta}}$ for all $f \in Y_{\theta}^{+}$. For, if $f \in Y_{\theta}^{+}$we can find $g \in X_{\theta}^{+}$with $\|g\|_{X_{\theta}}=1$ and $\left\langle f^{1 / 2}, g^{1 / 2}\right\rangle \geq(1-\varepsilon)\left\|f^{1 / 2}\right\|_{Y_{\theta}^{(2)}}=(1-$ ع) $\|f\|_{Y_{\theta}}^{1 / 2}$. We have

$$
\begin{aligned}
\left\langle g^{1 / 2}, g^{1 / 2}\right\rangle & \leq\left\langle(T f)^{1 / 2},(S g)^{1 / 2}\right\rangle \leq\left\|(T f)^{1 / 2}\right\|_{\bar{Y}_{\theta}^{(2)}}\left\|(S g)^{1 / 2}\right\|_{\bar{X}_{\theta}^{(2)}} \\
& \leq\left\|(T f)^{1 / 2}\right\|_{Z_{\theta}^{(2)}}\left\|(S g)^{1 / 2}\right\|_{\bar{X}_{\theta}^{(2)}}=\|T f\|_{Z_{\theta}}^{1 / 2}\|S g\|_{\bar{X}_{\theta}}^{1 / 2} \\
& \leq\|T f\|_{Z_{\theta}}^{1 / 2}\|S\|_{X_{\theta} \rightarrow \bar{X}_{\theta}}^{1 / 2}\|g\|_{X_{\theta}}^{1 / 2} ;
\end{aligned}
$$

hence $\|T f\|_{Z_{\theta}} \geq(1-\varepsilon)^{2}\|S\|_{X_{\theta} \rightarrow X_{\theta}}^{-1}\|f\|_{Y_{\theta}}$.
We thus have $\|T f\|_{Z_{\theta}} \simeq\|f\|_{Y_{\theta}}$ for all $f \in Y_{\theta}^{+}$(but $T$ need not to be an isomorphism for the metric structures of $Y_{\theta}^{+}$and $Z_{\theta}^{+}$).
(B) We are now in position to apply either Lemma 21 to $T: Y \rightarrow Z$ or Lemma 22 to $Z_{\theta}$. We suppose first that $Z \neq L_{\infty}$. Then the reasoning is very similar to that of [JMST], Thm. 5.6. For $n \geq 1$ and $1 \leq i \leq 2^{n}$ set $y_{n, i}=$ $\mathbf{1}_{\left[(i-1) 2^{-n}, i 2^{-n}\right]}$ and $z_{n}=\sup _{i=1}^{2^{n}} T y_{n, i}$. Note that $z_{n} \leq T 1$.

Case I. There exists $R$ such that $\inf _{n}\left\|\mathbf{1}_{\left\{z_{n} \leq R\right\}} z_{n}\right\|_{Z>0}$. Then by Lemma 21, there exists $K$ such that $\forall f \in Z_{+},\|f\|_{Z} \leq K\|f\|_{Y}$; the converse inequality holds by hypothesis, hence $Y=Z$.

Case II. There exist sequences $n_{l} \rightarrow \infty$ and $R_{l} \rightarrow \infty$ such that $\left\|\mathbf{1}_{\left\{z_{n_{l}} \leq R_{l}\right\}} z_{n_{l}}\right\|_{Z^{\prime \prime}} \rightarrow 0$ on $l \rightarrow \infty$. Then, after extracting if necessary, $z_{n_{l}} \rightarrow 0$ a.e. as $l \rightarrow \infty$. Since $z_{n_{l}} \leq T 1 \in L_{1}(1 / \psi)$, we have $\lim _{l \rightarrow \infty}\left\|z_{n_{l}}\right\|_{L_{1}(1 / \psi)}=0$ by Lebesgue's Theorem.

Then, by the interpolation inequality, $\left\|z_{n_{l}}\right\|_{z_{\theta}} \rightarrow 0$ as $l \rightarrow \infty$. If $y=$ $\sum_{i=1}^{2^{n}} b_{i} y_{n, i}$ is a simple nonnegative dyadic function, the norm of $T y$ in $Z_{\theta}$ is estimated by applying Lemma 22 to the $1 / \theta$-concave lattice $Z_{\theta}$ and to the finite sequence $\left(T y_{n, i}\right)_{1 \leq i \leq 2^{n}}$ (which is positively equivalent to $\left(y_{n, i}\right)_{1 \leq i \leq 2^{n}}$, hence positively symmetric), we obtain

$$
\|y\|_{Y_{\theta}} \leq C\left[\|y\|_{\infty}\left\|z_{n}\right\|_{z_{\theta}} \vee\|y\|_{L_{1}(00,1)}\right]
$$

where the constant does not depend on $n$. For a fixed dyadic function $y$, we make $n=n_{l}$ and let $l \rightarrow \infty$, deducing that $\|y\|_{Y_{\theta}} \leq C\|y\|_{1}$, which implies that $Y_{\theta}=L_{1}([0,1])$ (algebraically). Dualizing, we obtain $L_{\infty}=Y_{\theta}^{\prime}=\left[L_{1}, Y\right]_{\theta}^{\prime}=$ [ $\left.L_{\infty}, Y^{\prime}\right]_{\theta}=Y^{\prime 1 / \theta}$ (the $1 / \theta$-convexification of $Y^{\prime}$ ); hence $Y^{\prime}=L_{\infty}$, thus $Y=L_{1}$.
(C) We left aside in the preceding the case $Z=L_{\infty}$. In this case we may replace in the hypotheses the space $Z$ by the interpolation space $[Y, Z]_{1 / 2}=$ $Y^{(2)}$. We conclude that either $Y=Y^{(2)}$, which implies that $Y=L_{\infty}$, in contradiction with the hypothesis, or $Y=L_{1}$.

Proof of Theorem 16. The "if" part of Thm. 16 results from [LT2, Prop. $2 \mathrm{~d} 4]$. Conversely, let $X$ be a $q$-concave ( $q<2$ ) r.i. space over $\Omega=[0,1]$, and suppose that $X\left(l_{2}\right)$ embeds into $X$ as a complemented subspace. Denote by $X_{1 / 2}, X_{1 / 2}^{\prime}$ the $1 / 2$-concavifications of $X$ and $X^{\prime}$; note that $X_{1 / 2}$ is a 1/2-convex, concave quasi-Banach space, while $X_{1 / 2}^{\prime}$ is convex (up to renorming we may suppose that the quasi-norm on $X_{1 / 2}^{\prime}$ is a norm). Let $Z$ be the r.i. Banach function space over [0,1] defined by $f \in Z$ iff $f \otimes G^{2} \in$ $X_{1 / 2}^{\prime}$, where, as before, $G$ is a normal gaussian variable defined on the probability space $(S, \Sigma, \mathbf{P})$, equipped with the norm $\|f\|_{Z}=\| f \otimes$ $G^{2} \|_{X_{1 / 2}^{\prime}(\Omega \times S)}$. As conditional expectation operators are defined on the r.i. space $X_{1 / 2}^{\prime}$, we see that $Z \subset X_{1 / 2}^{\prime}$ (with $\|f\|_{Z} \geq(\mathbf{E}|G|) \cdot\|f\|_{X_{1 / 2}^{\prime}}$ ). We shall construct two bounded, positively linear operators $S: X_{1 / 2}^{+} \rightarrow X_{1 / 2}^{+}$and $T$ : $X_{1 / 2}^{\prime+} \rightarrow Z^{+}$verifying the "duality inequality"

$$
\forall f \in X_{1 / 2}^{+}, \forall g \in X_{1 / 2}^{+},\left\langle(T f)^{1 / 2},(S g)^{1 / 2}\right\rangle \geq\left\langle f^{1 / 2}, g^{1 / 2}\right\rangle
$$

(since $T$ is order continuous). An application of Prop. 18 then shows that $Z=X_{1 / 2}^{\prime}$; hence $X^{\prime}=Z^{(2)}$, and, by Lemma 17, $X^{\prime}$ has nontrivial upper Boyd index. By [LT2], Prop. 2b2, $X$ has a non-trivial lower Boyd index.
(A) Construction of $S$. Let

$$
X\left(l_{2}\right) \xrightarrow{V} X \xrightarrow{\pi} X\left(l_{2}\right)
$$

be a $C$-representation of $X\left(l_{2}\right)$ as a complemented subspace of $X$. For all $n$, we denote by $\left(x_{i, n}\right)_{i=1, \ldots 2^{n}}$ the dyadic partition of $[0,1]$ of order $n, x_{i, n}=$ $\left[(i=1) 2^{-n}, i 2^{-n}\right]$, and by $\left(e_{m}\right)$ the natural basis of $l_{2}$. Let $X_{n}$ be the sublattice of $X(\Omega)$ generated by the $x_{i, n}, i=1, \ldots 2^{n}$.

Let us first remark that in $X$, the unit ball of any hilbertian subspace $H$ is necessarily $X$-equiintegrable. For, if not, there exist a norm one sequence $\left(y_{k}\right)$ in $H$, a disjoint sequence $\left(y_{k}^{\prime}\right)$ in $X$, and a real $\delta>0$, with $\left|y_{k}^{\prime}\right| \leq\left|y_{k}\right|$ and $\left\|y_{k}^{\prime}\right\| \geq \delta$. Let $A_{k}$ be the support of $y_{k}^{\prime}$. Since $H$ is reflexive, we may suppose (up to extracting) that

$$
y_{k} \xrightarrow[k \rightarrow \infty]{w} y_{\infty}(\text { weakly })
$$

then

$$
z_{k}=y_{k}-y_{\infty} \xrightarrow{w} 0,
$$

and $z_{k}^{\prime}=\mathbf{1}_{A_{k}} z_{k}$ verifies $\lim \inf _{k \rightarrow \infty}\left\|z_{k}^{\prime}\right\|_{X} \geq \delta$; so we can suppose w.l.o.g.

$$
y_{k} \xrightarrow[k \rightarrow \infty]{w} 0
$$

hence (after extraction) that $\left(y_{k}\right)$ is equivalent to the $l_{2}$-basis. Then

$$
\begin{aligned}
\left\|\sum \alpha_{k} y_{k}^{\prime}\right\| & =\left\|\mathbf{E}_{\varepsilon}\left|\sum \alpha_{k} \varepsilon_{k} y_{k}^{\prime}\right|\right\| \leq\left\|\mathbf{E}_{\varepsilon}\left|\sum \alpha_{k} \varepsilon_{k} y_{k}\right|\right\| \\
& \leq \mathbf{E}_{\varepsilon}\left\|\sum \alpha_{k} \varepsilon_{k} y_{k}\right\| \sim\left(\sum\left|\alpha_{k}\right|^{2}\right)^{1 / 2}
\end{aligned}
$$

but by $q$-concavity of $X$,

$$
\left\|\sum \alpha_{k} y_{k}^{\prime}\right\|_{X} \geq c_{q}\left(\sum\left\|\alpha_{k} y_{k}^{\prime}\right\|_{X}^{q}\right)^{1 / q} \geq c_{q} \delta\left(\sum\left|\alpha_{k}\right|^{q}\right)^{1 / q}
$$

which is a contradiction since $q<2$.
For fixed $i$ and $n$, the sequence $\left(V\left(x_{i, n} \otimes e_{m}\right)\right)_{m=1}^{\infty}$ generates a subspace $E_{i, n}$ of $X$ which is complemented and isomorphic to $l_{2}$. By the preceding, the unit ball of $E_{i, n}$ is $X$-equiintegrable, and in particular the $L_{1}$-norm and the $X$-norm are equivalent on this subspace. Thus there exist in the unit ball of $E_{i, n}$ a sequence $\left(u_{i, n}^{m}\right)_{m}$ which converges w.c.d. to a (non-zero) conditionally Gaussian variable $A_{i, n} \otimes G$. Doing the same for each $i=1, \ldots 2^{n}$, we find $2^{n}$ sequences $\left(u_{i, n}^{m}\right)_{m}$, each converging w.c.d. to a variable $A_{i, n} \otimes G$. Hence we have the following joint w.c.d. convergence:

$$
\left(u_{i, n}^{m_{i}}\right)_{i=1}^{2^{n}} \xrightarrow[m_{1} \rightarrow \infty ; m_{2} \rightarrow \infty ; \ldots m_{2^{n} \rightarrow \infty}]{\mathrm{wcd}}\left(A_{i, n} \otimes G_{i}\right)_{i=1}^{2^{n}}
$$

where now the $G_{i}$ are independent normal gaussian variables. Reasoning as in $\S 3$, we obtain $2^{n}$ sequences $\left(v_{i, n}^{m}\right)_{m \geq 1}, i=1, \ldots 2^{n}$, in $X(\Omega \times S)$ which are jointly equimeasurable with the sequences $\left(A_{i, n} \otimes G_{i}^{m}\right)_{m \geq 1}$ and generating a $C$-representation of $X_{n}\left(l_{2}\right)$ as a complemented subspace of $X(\Omega \times S)$. Again an argument using a measure preserving transform of $\Omega \times S$ allows to replace the $\left(v_{i, n}^{m}\right)_{m \geq 1,1 \leq i \leq 2^{n}}$ by the $\left(A_{i, n} \otimes G_{i}^{m}\right)$. We have thus a factorization of the identity of $X_{n}\left(l_{2}\right)$ :

$$
X_{n}\left(l_{2}\right) \xrightarrow{j_{n}} X(\Omega \times S) \xrightarrow{\pi_{n}} X_{n}\left(l_{2}\right)
$$

with $j_{n}\left(x_{i, n} \otimes e_{m}\right)=A_{i, n} \otimes G_{i}^{m}$ and $\left\|j_{n}\right\|\left\|\pi_{n}\right\| \leq C$.
There is no relation between the systems $\left(A_{i n}\right)_{i=1, \ldots 2^{n}}$ for different values of $n$. However $j_{n}$ induces naturally a $C$-representation of $X_{n-1}\left(l_{2}\right)$ in $X(\Omega \times$ $S$ ),

$$
X_{n-1}\left(l_{2}\right) \xrightarrow{k_{n}} X(\Omega \times S) \xrightarrow{\sigma_{n}} X_{n-1}\left(l_{2}\right),
$$

with $k_{n}=j_{n} \circ i_{n}, \sigma_{n}=q_{n} \circ \pi_{n}$, where $i_{n}$ is the natural injection of $X_{n-1}\left(l_{2}\right)$ into $X_{n}\left(l_{2}\right)$, and $q_{n}$ the expectation projection from $X_{n}\left(l_{2}\right)$ onto $X_{n-1}\left(l_{2}\right)$.

Then

$$
\begin{aligned}
k_{n}\left(x_{i, n-1} \otimes e_{m}\right) & =j_{n}\left(\left(x_{2 i-1, n}+x_{2 i, n}\right) \otimes e_{m}\right) \\
& =A_{2 i-1, n} \otimes G_{2 i-1}^{m}+A_{2 i, n} \otimes G_{2 i}^{m}
\end{aligned}
$$

Hence the sequence $\left(k_{n}\left(x_{i, n-1} \otimes e_{m}\right)_{m}, i=1, \ldots 2^{n}\right.$, is jointly equimeasurable with $\left(A_{i, n-1}^{(n)} \otimes G_{i}^{m}\right)_{m}$, where $A_{i, n-1}^{(n)}=\left(A_{2 i-1, n}^{2}+A_{2 i, n}^{2}\right)^{1 / 2}$. Again (with a measure-preserving transform on $\Omega \times S$ ), we deduce another $C$-representation of $X_{n-1}\left(l_{2}\right)$ given by

$$
X_{n-1}\left(l_{2}\right) \xrightarrow{j_{n}^{(n)}} X(\Omega \times S) \xrightarrow{\pi_{n}^{(n)}} X_{n}\left(l_{2}\right)
$$

with $j_{n-1}^{(n)}\left(x_{i, n-1} \otimes e_{m}\right)=A_{i, n-1}^{(n)} \otimes G_{i}^{m}\left(m \geq 1, i=1, \ldots 2^{n-1}\right)$. In the same way we recursively define $j_{k}^{(n)}, \pi_{k}^{(n)}$, for each $k \leq n$, giving a $C$-representation of $X_{k}\left(l_{2}\right)$ as a complemented subspace of $X(\Omega \times S)$.

We claim that for fixed $k, i, m$, the sequence $\left(j_{k}^{(n)}\left(x_{i, k} \otimes e_{m}\right)\right)_{n \geq k}=$ $\left(A_{i, k}^{(n)} \otimes G_{i}^{m}\right)_{n \geq k}$ is $X$-equiintegrable. For each $n, A_{i, k}^{(n)} \otimes G_{i}^{m}$ is equimeasurable with

$$
u_{i, k}^{(n)}=\sum_{(i-1) 2^{-k} \leq j 2^{-n} \leq i 2^{-k}} A_{j, n} \otimes G_{j}^{n}
$$

which by construction is the wcd limit of a sequence $\left(V\left(f_{m}^{n}\right)\right)_{m \geq 1}$, where the elements $f_{m}^{n}$ are disjoint in the lattice $X\left(l_{2}\right)$ but verify $\left\|f_{m}^{n}\right\|_{2}(\omega)=x_{i, k}(\omega)$ a.e. In fact we may even suppose that $f_{m}^{n} \in X\left(H_{m}^{n}\right)$, where $H_{m}^{n}=\overline{\operatorname{span}}\left[e_{l}\right]_{l \in L_{m}^{n}}$ and $L_{1}^{n}<L_{2}^{n}<\cdots$ are successive disjoint intervals of $\mathbf{N}$. If the sequence $\left(u_{i, k}^{(n)}\right)_{n}$ is not $X$-equiintegrable, it is possible to find sequences $n_{1}<\cdots<$ $n_{l}<n_{l+1}<\cdots$ and $m_{1}<\cdots<m_{l}<m_{l+1}<\cdots$ such that the sequence $\left(V\left(f_{m_{l}}^{n_{l}}\right)\right)_{l}$ is not $X$-equiintegrable, while $\left(f_{m_{l}}^{n_{l}}\right)_{l}$ is disjoint in the lattice $X\left(l_{2}\right)$. But then

$$
\left\|\sum_{l} \alpha_{l} f_{m_{l}}^{n_{l}}\right\|_{2}(\omega)=\left(\sum_{l} \alpha_{l}^{2}\right)^{1 / 2} x_{i, k}(\omega) \quad \text { a.e. }
$$

hence $\left(f_{m_{l}}^{n_{l}}\right)_{l}$ is equivalent in $X$ to the $l_{2}$-basis; so is $\left(V\left(f_{m_{l}}^{n_{l}}\right)_{l}\right.$, which is then $X$-equiintegrable by a preceding remark, a contradiction which proves our claim.

The sequence $\left(j_{k}^{(n)}\left(x_{i, k} \otimes e_{m}\right)\right)_{n \geq k}$ defines an element $\tilde{J}_{k}\left(x_{i, k} \otimes e_{m}\right)$ of the ultrapower $\hat{X}=X(\Omega \times S)^{\mathrm{N}} / \mathscr{U}$, in fact of $\hat{X}^{\text {eq }}$. Similarly the sequence $\left(A_{i, k}^{(n)}\right)_{n \geq k}$ (which is a fortiori $X$-equiintegrable) defines an element $\tilde{A}_{i, k}$ of $\tilde{X} \xlongequal{=} X(\Omega)^{\mathrm{N}} / \mathscr{U}$, in fact of $\tilde{X}^{\text {eq }}$. We identify as usual $\tilde{X}^{\text {eq }}=X(\tilde{\Omega}, \tilde{\mathscr{A}}, \tilde{\mu})$. We have a natural isometric embedding $\Theta: X(\tilde{\Omega} \times S) \rightarrow \hat{X}^{\mathrm{eq}}:$ if $\tilde{f} \in \tilde{X}^{\mathrm{eq}}$ is
defined by the sequence $\left(f_{n}\right)_{n}$ and $g \in L_{\infty}(S)$, then $\Theta(\tilde{f} \otimes g)$ is simply the element of $\hat{X}^{\text {eq }}$ defined by the sequence $\left(f_{n} \otimes g\right)_{n}$. We have clearly $\Theta\left(\tilde{A}_{i, k} \otimes G_{i}^{m}\right)=\tilde{J}_{k}\left(x_{i, k} \otimes e_{m}\right)$.

Now define $\tilde{\pi}_{k}: \hat{X} \rightarrow X_{k}\left(l_{2}\right)$ by $\tilde{\pi}_{k}(\tilde{h})=w-\lim _{n, \mathscr{U}} \pi_{k}^{(n)}\left(h_{n}\right)$ when $\tilde{h}$ is defined by the sequence $\left(h_{n}\right)_{n}$. We now have a $C$-factorization of the identity of $X_{k}\left(l_{2}\right)$ through $\hat{X}^{\mathrm{eq}}$ :

$$
X_{k}\left(l_{2}\right) \xrightarrow{\tilde{j}_{k}} \hat{X}^{\mathrm{eq}} \xrightarrow{\tilde{\pi}_{k}} X_{k}\left(l_{2}\right)
$$

In fact we can replace $\hat{X}^{\text {eq }}$ by $X(\tilde{\Omega} \times S)$, replacing $\tilde{J}_{k}$ by $\Theta^{-1} \circ \tilde{J}_{k}$ and $\tilde{\pi}_{k}$ by $\tilde{\pi}_{k} \circ \Theta$. But from

$$
A_{i, k}^{(n)}=\left(A_{2 i-1, k+1}^{(n) 2}+A_{2 i, k+1}^{(n) 2}\right)^{1 / 2}, \quad \forall n \geq k+1
$$

we get

$$
\tilde{A}_{i, k}=\left(\tilde{A}_{2 i-1, k+1}^{2}+\tilde{A}_{2 i, k}^{2}\right)^{1 / 2}
$$

Since these $\tilde{A_{i, k}}$ live in a separable sublattice of $X(\tilde{\Omega})$, we can find a new family $\left(\bar{A}_{i, k}\right)_{k \geq 1 ; i=1, \ldots 2^{k}}$ in $X(\Omega)$, jointly equimeasurable with $\left(\tilde{A}_{i, k}\right)_{i, k}$, and sharing the same properties. So (writing $A_{i, k}$ in place of $\bar{A}_{i, k}$ ) we are back to our starting point, i.e., we have a family $\left(A_{i, k}\right)_{i, k}$ such that $\left(A_{i, k} \otimes\right.$ $\left.G_{i}^{m}\right)_{m \geq 1,1 \leq i \leq 2^{k}}$ generates a complemented subspace of $X(\Omega \times S)$ isomorphic to $X_{k}\left(l_{2}\right)$, but we gained the compatibility condition

$$
\forall k, \forall i, 1 \leq i \leq 2^{k}, \quad A_{i, k}^{2}=A_{2 i-1, k+1}^{2}+A_{2 i, k+1}^{2}
$$

We define the operator $S$ on positive dyadic functions by

$$
f=\sum_{i=1}^{2^{k}} \lambda_{i}^{2} x_{i, k} \Rightarrow S(f)=\sum_{i=1}^{2^{k}} \lambda_{i}^{2} A_{i, k}^{2}
$$

Due to the compatibility condition, $S(f)$ does not depend on the way of writing $f$ as a simple dyadic function. We have

$$
\begin{aligned}
\left\|S\left(\sum_{i} \lambda_{i}^{2} x_{i, k}\right)\right\|_{X_{1 / 2}} & =\left\|\left(\sum_{i} \lambda_{i}^{2} A_{i, k}^{2}\right)^{1 / 2}\right\|_{X}^{2} \sim\left\|\left(\sum_{i} \lambda_{i}^{2} A_{i, k}^{2}\right)^{1 / 2} \otimes G\right\|_{X}^{2} \\
& =\left\|\sum_{i} \lambda_{i} A_{i, k} \otimes G_{i}^{k}\right\|_{X}^{2} \sim\left\|\sum_{i} \lambda_{i} x_{i, k}\right\|_{X}^{2}=\left\|\sum_{i} \lambda_{i}^{2} x_{i, k}\right\|_{X_{1 / 2}}
\end{aligned}
$$

thus $\|S f\|_{X_{1 / 2}} \sim\|f\|_{X_{1 / 2}}$ for every positive dyadic function $f$. Note that $|S f-S g| \leq S(f \vee g)-S(f \wedge g)=S(|f-g|)$, by positivity of $S$; thus $S$ is lipschitzian on the cone of positive dyadic functions, and it extends by density
to a positively linear operator $X_{1 / 2} \rightarrow X_{1 / 2}$, with again $\|S f\|_{X_{1 / 2}} \sim\|f\|_{X_{1 / 2}}$ for all $f$.
(B) Construction of $T$. For each $k \geq 1$, we have a projection

$$
P_{k}: X(\Omega \times S) \rightarrow E_{k}=\overline{\operatorname{span}}\left[A_{i, k} \otimes G_{i}^{m}\right]_{m \geq 1, i=1, \ldots 2^{k}}
$$

From now on, we suppose that $(\Sigma, \sigma)$ is generated by the variables $\left(G_{i}^{m}\right)_{i, m}$. We consider the action of the group $G_{n, p}=O(n)^{p}$ on the space $X(\Omega \times S)$, defined by

$$
\begin{aligned}
& \left(U_{1}, \ldots U_{p}\right) f\left(\omega,\left(G_{1}^{m}\right)_{m},\left(G_{2}^{m}\right)_{m} \ldots\left(G_{p}^{m}\right)_{m},\left(G_{p+1}^{m}\right)_{m}, \ldots\right) \\
& \quad=f\left(\omega, U_{1}^{*}\left(G_{1}^{m}\right)_{m}, U_{2}^{*}\left(G_{2}^{m}\right)_{m}, \ldots U_{p}^{*}\left(G_{p}^{m}\right)_{m},\left(G_{p+1}^{m}\right)_{m}, \ldots\right)
\end{aligned}
$$

where, as in $\S 2, O(n)$ acts on $\mathbf{R}^{\mathbf{N}}$ by $U\left(x_{1}, \ldots, x_{n}, x_{n+1} \ldots\right)=$ $\left(U\left(x_{1}, \ldots x_{n}\right), x_{n+1} \ldots\right)$. The subspace $E_{k}$ is invariant under the action of $G_{n, p}$ (in fact each "fiber" $E_{k, i}=\overline{\operatorname{span}}\left[A_{i, k} \otimes G_{i}^{m}\right]_{m \geq 1}$ is invariant). The reasoning of Proposition 6 gives us a projection $R_{k}: X(\Omega \times S) \rightarrow E_{k}$ which is invariant under the action of all the groups $G_{n, p}$; the method of Proposition 6 shows that $R_{k}$ takes necessarily the form

$$
R_{k} f=2^{k} \sum_{\substack{1 \leq i \leq 2^{k} \\ m \geq 1}}\left\langle f, B_{i, k} \otimes G_{i}^{m}\right\rangle A_{i, k} \otimes G_{i}^{m}
$$

with $B_{i, k} \otimes G_{i}^{m} \in X^{\prime}$ and $\left\langle A_{i, k}, B_{i, k}\right\rangle=2^{-k}$.
Since $\left(B_{i, k} \otimes G_{i}^{m}\right)_{i, m}$ is biorthogonal to $\left(B_{i, k} \otimes G_{i}^{m}\right)_{i, m}$, which is equivalent to $\left(x_{i, k} \otimes e_{m}\right)_{i, m}$ and spans a complemented subspace of $X(\Omega \times S)$, and since $\left\langle B_{i, k}, A_{i, k}\right\rangle \stackrel{m}{=} 2^{-k}=\left\langle x_{i, k}, x_{i, k}\right\rangle$, we see that

$$
\left\|\sum_{i, m} \alpha_{i, m} B_{i, k} \otimes G_{i}^{m}\right\|_{X^{\prime}} \sim\left\|\sum_{i, m} \alpha_{i, m} x_{i, k} \otimes e_{m}\right\|_{X^{\prime}\left(l_{2}\right)}
$$

Again there is a priori no relation between the systems $\left(B_{i, k}\right)_{i=1, \ldots 2^{k}}$ for different values of $k$, so we modify them to have a compatibility condition. The method we used for the ( $A_{i, k}$ ) does not work here because $X^{\prime}$ has no more non trivial concavity. However we set

$$
B_{i, k-1}^{(k)}=\left(B_{2 i-1, k}^{2}+B_{2 i, k}^{2}\right)^{1 / 2}
$$

and note that $\left(B_{i, k}^{(k)} \otimes G_{i}^{m}\right)_{m \geq 1,1 \leq i \leq 2^{k-1}}$ span an isomorph of $X_{k-1}^{\prime}\left(l_{2}\right)$ in $X^{\prime}$ (with no loss on the equivalence constant), and that

$$
\begin{aligned}
\left\langle B_{i, k-1}^{(k)}, A_{i, k-1}\right\rangle & =\left\langle\left(B_{2 i-1, k}^{2}+B_{2 i, k}^{2}\right)^{1 / 2},\left(A_{2 i-1, k}^{2}+A_{2 i, k}^{2}\right)^{1 / 2}\right\rangle \\
& \geq\left\langle B_{2 i-1, k}, A_{2 i-1, k}\right\rangle+\left\langle B_{2 i, k}, A_{2 i, k}\right\rangle=2^{-(k-1)}
\end{aligned}
$$

Similarly we recursively define $B_{i, p}^{(k)}, p \leq k$, which satisfies

$$
\left\|\sum_{\substack{1 \leq i \leq 2^{p} \\ m \geq 1}} \alpha_{i, m} B_{i, p}^{(k)} \otimes G_{i}^{m}\right\|_{X^{\prime}} \sim\left\|\sum_{\substack{1 \leq i \leq 2^{p} \\ m \geq 1}} \alpha_{i, m} x_{i, p} \otimes e_{m}\right\|_{X^{\prime}}
$$

(with equivalence constants independent from $k, p$ ) and

$$
\left\langle B_{i, p}^{(k)}, A_{i, p}\right\rangle \geq 2^{-p} .
$$

The sequence $\left(B_{i, p}^{(k) 2}\right)$ is bounded in the lattice $X_{1 / 2}^{\prime}$, which is a dual lattice (since it is $r$-convex, $r>1$, and a maximal r.i. space, $X_{1 / 2}^{\prime}=Z^{*}$ with $\left.Z=\left(X_{1 / 2}^{\prime}\right)^{\prime}\right)$. We define $\hat{B}_{i, p}^{2}=w^{*}-\lim _{k, \mathscr{U}} B_{i, p}^{(k) 2}$. We have $\hat{B}_{i, p} \in X^{\prime}$ and in fact $\hat{B}_{i, p} \otimes G \in X^{\prime}$, with $\hat{B}_{i, p}^{2} \otimes G^{2}=w^{*}-\lim _{k, \mathscr{U}} B_{i, p}^{(k) 2} \otimes G^{2}$. Thus

$$
\begin{aligned}
\left\|\sum_{i, m} \alpha_{i, m} \hat{B}_{i, p} \otimes G_{i}^{m}\right\|_{X^{\prime}} & =\left\|\left(\sum_{i, m} \alpha_{i, m}^{2} \hat{B}_{i, p}^{2}\right)^{1 / 2} \otimes G\right\|_{X^{\prime}} \\
& =\left\|\sum_{i, m} \alpha_{i, m}^{2} \hat{B}_{i, p}^{2} \otimes G^{2}\right\|_{X_{1 / 2}^{\prime}}^{1 / 2} \\
& \leq \lim _{k, \mathscr{U}}\left\|\sum_{i, m} \alpha_{i, m}^{2} B_{i, p}^{(k) 2} \otimes G^{2}\right\|_{X_{1 / 2}^{\prime}}^{1 / 2} \\
& =\lim _{k, \mathscr{U}}\left\|\sum_{i, m} \alpha_{i, m} B_{i, p}^{(k)} \otimes G_{i}^{m}\right\|_{X^{\prime}} \\
& \sim\left\|\sum_{i, m} \alpha_{i, m} x_{i, p} \otimes e_{m}\right\|_{X^{\prime}\left(l_{2}\right)}
\end{aligned}
$$

Set $C_{i, p}=w^{*}-\lim _{k, \mathscr{U}} B_{i, p}^{(k)}$, where now the $w^{*}$-limit is relative to $\sigma\left(X^{\prime}, X\right)$.

We have $C_{i, p} \leq \hat{B}_{i, p}$, as is well known; hence, for nonnegative reals $\alpha_{i, m}$,

$$
\begin{aligned}
& \left\langle\sum_{i, m} \alpha_{i, m} \hat{B}_{i, p} \otimes G_{i, m}, \sum_{j, l} \beta_{j, l} A_{j, p} \otimes G_{j}^{l}\right\rangle \\
& \quad=\sum_{i, m} \alpha_{i, m} \beta_{i, m}\left\langle\hat{B}_{i, p}, A_{i, p}\right\rangle \\
& \quad \geq \sum_{i, m} \alpha_{i, m} \beta_{i, m}\left\langle C_{i, p}, A_{i, p}\right\rangle \\
& \quad=\lim _{k, \mathscr{U}_{i, m}} \sum_{i, m} \alpha_{i, m}\left\langle B_{i, p}^{(k)}, A_{I, p}\right\rangle \geq 2^{-p} \sum_{i, m} \alpha_{i, m} \beta_{i, m} \\
& \quad=\left\langle\sum_{i, m} \alpha_{i, m} x_{i, p} \otimes e_{m}, \sum_{j, l} \beta_{j, l} x_{j, p} \otimes e_{l}\right\rangle
\end{aligned}
$$

And since

$$
\left\|\sum_{j, l} \beta_{j, l} A_{j, p} \otimes G_{j}^{l}\right\|_{X} \sim\left\|\sum_{j, l} \beta_{j, l} x_{j, p} \otimes e_{l}\right\|_{X\left(l_{2}\right)}
$$

we have

$$
\left\|\sum_{i, m} \alpha_{i, m} \hat{B}_{i, p} \otimes G_{i, m}\right\|_{X^{\prime}} \gtrsim\left\|\sum_{i, m} \alpha_{i, m} x_{i, p} \otimes e_{m}\right\|_{X^{\prime}\left(l_{2}\right)}
$$

Since we have, by construction, $B_{i, p}^{(k) 2}=B_{2 i-1, p+1}^{(k) 2}+B_{2 i, p+1}^{(k) 2}$, for $k \geq p+1$, we obtain, by passing to the limit, the compatibility condition

$$
\hat{B}_{i, p}^{2}=\hat{B}_{2 i-1, p+1}^{2}+\hat{B}_{2 i, p+1}^{2}
$$

We can now define $T f$ for a nonnegative dyadic function $f=\sum_{i=1}^{2^{p}} \alpha_{i}^{2} x_{i, p}$ by $T f=\sum_{i=1}^{2^{p}} \alpha_{i}^{2} \hat{B}_{i, p}^{2}$, and we have

$$
\left\|(T f) \otimes G^{2}\right\|_{X_{1 / 2}^{\prime}} \sim\|f\|_{X_{1 / 2}^{\prime}}
$$

and

$$
\left\langle(T f)^{1 / 2},(S g)^{1 / 2}\right\rangle \geq\left\langle f^{1 / 2}, g^{1 / 2}\right\rangle
$$

for every nonnegative dyadic $f$ and $g$.
Let $Z$ be the space $\left\{f \in L_{0} \mid f \otimes G \in X_{1 / 2}^{\prime}\right\}$. We can extend (by density) the operator $T$ to an operator $E_{+} \rightarrow Z$, where $E$ is the closure of dyadic functions in $X_{1 / 2}^{\prime}$ which is also the closure of $L_{\infty}$ in $X_{1 / 2}^{\prime}$ (since $X_{1 / 2}^{\prime} \neq L_{\infty}$, i.e. $X \neq L_{1}$ ). If $f \in X_{1 / 2}^{\prime}$ and $f_{n} \in E_{+}, f_{n} \uparrow f$ a.e., we have

$$
\forall n,\left\|T f_{n}\right\|_{Z} \leq C\left\|f_{n}\right\|_{X_{1 / 2}^{\prime}} \leq C\|f\|_{X_{1 / 2}^{\prime}}
$$

hence $F=\sup _{n} T f_{n}$ belongs to $Z^{\prime \prime}$, with norm $\|F\|_{Z^{\prime \prime}} \leq C\|f\|_{X_{1 / 2}^{\prime}}$. This element $F$ does not depend on the sequence $f_{n} \uparrow f$, since if $0 \leq g_{m} \uparrow f$ we have: $g_{m} \wedge f_{n} \uparrow f_{n}$ a.e. and in the norm of $E$, hence $T\left(g_{m} \wedge f_{n}\right) \rightarrow T f_{n}$ as $m \rightarrow \infty$ and $\sup _{m} T g_{m} \geq \sup _{m, n} T\left(g_{m} \wedge f_{n}\right) \geq \sup _{n} T f_{n}$. We set $T f=F$. The operator $T$ is then positively additive and order continuous, and, choosing $\left(f_{n}\right)$ such that $\left\|f_{n}\right\| \rightarrow\|f\|$, we have

$$
\|T f\|_{Z} \geq \sup _{n}\left\|T f_{n}\right\|_{Z} \sim \sup _{n}\left\|f_{n}\right\|_{X_{1 / 2}^{\prime}}=\|f\|_{X_{1 / 2}^{\prime}} .
$$

This ends the proof of Theorem 16.
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