

Operators Between Subspaces and Quotients of L^1

G. GODEFROY, N. J. KALTON & D. LI

ABSTRACT. We provide an unified approach of results of L. Dor on the complementation of the range, and of D. Alspach on the nearness from isometries, of small into-isomorphisms of L^1 . We introduce the notion of small subspace of L^1 , and show lifting theorems for operators between quotients of L^1 by small subspaces. We construct a subspace of L^1 which shows that extension of isometries from subspaces of L^1 to the whole space are no longer true for isomorphisms, and that nearly isometric isomorphisms from subspaces of L^1 into L^1 need not be near from any isometry.

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I. Introduction. The Banach space L^1 is of fundamental importance for Fourier analysis, probability theory, and many other fields of pure and applied analysis. However, its Banach space structure is not yet fully understood. In particular, the study of operators from L^1 to L^1 lead to simply stated but apparently very hard problems: for instance, it is not known whether a complemented subspace of L^1 is necessarily isomorphic to L^1 or ℓ^1 (see [35]). A major tool for handling operators on $L^1(\Omega)$ is their representation through a “matrix” whose “rows” are indexed by the points of Ω and consist of measures on Ω ([18]; see also [9]). This representation

will play a major role in the present work, together with duality arguments, which will allow us to use methods and results from the isometric theory of Banach spaces. We will focus on several natural aspects of operator theory on L^1 : perturbation of operators which are close to isometries, lifting of operators between quotient spaces, extensions of operators defined on linear subspaces. Such questions have of course been considered before and important theorems have been obtained (see e.g. [1], [7], [8], [15], [19], [21], [32], and references therein). Our work provides on one hand alternative simpler proofs to some known theorems, which are subsequently improved, and on the other hand new statements. We now turn to a detailed description of our results.

Section II begins with a simple inequality, whose proof relies on Gaussian randomization, which provides a control of the atomic part of norm-increasing operators (Lemma II.1). This inequality leads to a unified approach to Dor's complementation theorem (Theorem II.3) and to Alspach's perturbation theorem (Theorem II.7). These two results follow from a Hahn-Banach argument applied to the atomic part of the relevant operator, considered as taking values in a vector-valued L^1 . Let us mention that our approach provides a quantitative improvement of Alspach's theorem, but do not improves the constant $\sqrt{2}$ in the statement of Dor's result. The new notion of "small subspace" is introduced in Section III (Definition III.1). Roughly speaking, a "small subspace" is a subspace of L^1 which is nowhere "locally equal" to L^1 . Operators which have a non trivial "diagonal" in their "matrix" representation are denoted strong Enflo operators (Definition III.3). The link between these two notions is that operators which take their values in a small subspace are not strong Enflo; in fact, more is true (Proposition III.6). Small subspaces are studied within the class of translation invariant subspaces; it is also shown that a closed direct sum of two small subspaces is small (Proposition III.9). Smallness of spaces is a crucial concept in Section IV, which is devoted to the lifting of operators between quotient spaces: the main result of this section is that, if X and Y are small subspaces, and the quotient spaces L^1/X and L^1/Y are close to each other in Banach-Mazur distance, then there is an isomorphism of L^1 which maps X close to Y with respect to the Hausdorff distance (Theorem IV.3). When one of the spaces X or Y has the lifting property (e.g. if one of these spaces is complemented in its bidual), it follows that there is an invertible operator of L^1 which maps X onto Y (Corollary IV.8). This result can be seen as a substitute in L^1 of a theorem of Lindenstrauss and Rosenthal on ℓ^1 ([24]; see [25], Th.2.f.8). Note, however, that the quantitative behaviour of the liftings is actually better in the L^1 case (see Remark IV.9). Finally,

Section V provides a counterexample, showing the “sensitivity over ε ” of the positive results obtained so far about isometries: there is a subspace X of L^1 such that, although isometries from this space into L^1 extend to isometries of the whole space L^1 , there are linear maps T from X into L^1 , which are arbitrarily close to isometries, and such that no small perturbation of T is actually an isometry from X into L^1 ; in particular, these “near isometries” do not extend to near isometries defined on L^1 (Theorem V.1). This space X is arbitrarily close in the Banach-Mazur sense to weak-star closed subspaces of ℓ^1 , though it is far in the Hausdorff distance of any subspace of L^1 spanned by disjoint functions; hence, although it has a strong “ ℓ^1 ” behaviour (for instance, its unit ball is compact locally convex in measure, see [12]), it cannot be decomposed into essentially disjoint finite dimensional parts (compare with [7], [33]). The construction of X makes a crucial use of p -stable random variables with p close to 1.

Our notation is classical. All the Banach spaces we consider are indifferently real or complex, except in Section V, where we work only in real case. All the measure spaces are separable and purely non-atomic, and the measures are (unless other stated) probability measures; hence the corresponding L^1 -spaces are isometric to $L^1(0, 1)$, but for convenience we use sometimes a more general probability space. When writing $L^1([0, 1], m)$, we understand that m is a Lebesgue measure. The following representation of operators on L^1 ([18], Theorem 3.1) will be very useful: for every operator $T : L^1 \rightarrow L^1$ there are measures μ_x , with $x \mapsto \mu_x$ weak*-measurable, such that, for almost every x ,

$$(R) \quad Tf(x) = \int_0^1 f(s) d\mu_x(s) = \sum_{n=1}^{+\infty} a_n(x) f(\sigma_n(x)) + \int_0^1 f(s) dv_x(s),$$

where $\sum_{n=1}^{+\infty} a_n(x) \delta_{\sigma_n(x)}$ and v_x are the atomic and continuous part of μ_x .

Section III and IV are essentially independent of Section II.

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II. Projections and small into-isomorphisms of L^1 . In this section, we provide a unified approach to results of D. Alspach and L. Dor on isomorphisms from L^1 into L^1 with small bounds. A short proof of the special case of convolutors, which points towards the main lemma, will be given at the end of the section (see Remark II.8). The key result of this approach is the following lemma:

Lemma II.1. *Let $T : L^1([0, 1], \nu) \rightarrow L^1([0, 1], \nu)$, and write*

$$Tf(s) = \sum_{n=1}^{\infty} a_n(s)f(\sigma_n(s)) + \int_0^1 f(t) d\nu_s(t).$$

Suppose that $\|Tf\|_1 \geq \|f\|_1$ for all $f \in L^1$. Then, for all $f \in L^1$,

$$\|f\|_1 \leq \left\| \left(\sum_{k=1}^{\infty} |a_k(s)|^2 |f(\sigma_k(s))|^2 \right)^{1/2} \right\|_1.$$

Proof. Let (y_k) be a sequence of i.i.d. Gaussian (real or complex) random variables on a probability space (Ω, \mathbb{P}) , normalized in L^1 . For any $l > 0$, we let

$$D_k^l = \left[\frac{k-1}{2^l}, \frac{k}{2^l} \right], \quad 1 \leq k \leq 2^l,$$

be the dyadic intervals of level l .

Fact. Let ν be a diffuse measure on $[0, 1]$. Defining

$$\varphi_l = \sum_{k=1}^{2^l} |\nu(D_k^l)|^2,$$

one has $\lim_l \varphi_l = 0$.

Indeed, $\varphi_l \leq \|\nu\|_1 \cdot \max_k (|\nu(D_k^l)|)$, since $(\sum |\lambda_k|^2) \leq (\sum |\lambda_k|) \cdot (\max |\lambda_k|)$ for every sequence (λ_k) of scalars. Since ν is diffuse, this \max_k tends to 0 when l goes to infinity, and the conclusion follows. \square

Let now ρ be an arbitrary measure on $[0, 1]$. Let us write

$$\rho = \sum_{n=1}^{\infty} a_n \delta_{x_n} + \nu = \alpha + \nu,$$

where α is atomic and ν is diffuse. For any l and $N = 2^l$, one has

$$\sum_{k=1}^N |\rho(D_k^l)|^2 = \sum_{k=1}^N |\alpha(D_k^l) + \nu(D_k^l)|^2,$$

and from the Fact, it follows easily that

$$(1) \quad \lim_N \sum_{k=1}^N |\rho(D_k^l)|^2 = \sum_{n=1}^{\infty} |a_n|^2.$$

We now denote $D_k^l = D_k$, and for a given $N = 2^l$, any $\omega \in \Omega$, and any $f \in L^1$, we define

$$f_\omega = \sum_{k=1}^N \gamma_k(\omega) \mathbf{1}_{D_k} f.$$

With this notation, one has

$$\begin{aligned} \int_{\Omega} \|f_\omega\|_1 d\mathbb{P}(\omega) &= \iint_{[0,1] \times \Omega} |f_\omega(t)| dt d\mathbb{P}(\omega) \\ &= \iint_{[0,1] \times \Omega} \left| \sum_{k=1}^N \gamma_k(\omega) \mathbf{1}_{D_k}(t) f(t) \right| dt d\mathbb{P}(\omega) \\ &= \int_{\Omega} \left(\sum_{k=1}^N |\gamma_k(\omega)| \int_{D_k} |f(t)| dt \right) d\mathbb{P}(\omega) = \|f\|_1, \end{aligned}$$

and it follows from the assumption on T that, for any f , one has

$$\|f\|_1 \leq \int_{\Omega} \|T(f_\omega)\|_1 d\mathbb{P}(\omega).$$

We now use the representation (R) of T

$$(R) \quad Tf(s) = \int_0^1 f(t) d\mu_s(t) = \sum_{n=1}^{\infty} a_n(s) f(\sigma_n(s)) + \int_0^1 f(t) d\nu_s(t),$$

and for any given $f \in L^1$ we define, for all s , a measure ρ_s by $d\rho_s = f \cdot d\mu_s$.

With this notation, one has

$$T(f_\omega)(s) = \sum_{k=1}^N \gamma_k(\omega) \rho_s(D_k),$$

and thus, by Fubini's Theorem,

$$\begin{aligned} \int_{\Omega} \|T(f_\omega)\|_1 d\mathbb{P}(\omega) &= \int_0^1 \int_{\Omega} \left| \sum_{k=1}^N \gamma_k(\omega) \rho_s(D_k) \right| d\mathbb{P}(\omega) ds \\ &= \int_0^1 \left(\sum_{k=1}^N |\rho_s(D_k)|^2 \right)^{1/2} ds. \end{aligned}$$

It follows that

$$\|f\|_1 \leq \int_0^1 \left(\sum_{k=1}^N |\rho_s(D_k)|^2 \right)^{1/2} ds.$$

On the other hand, by applying (1) to $\rho = \rho_s$, we have for almost every s :

$$\lim_N \left(\sum_{k=1}^N |\rho(D_k)|^2 \right)^{1/2} = \left(\sum_{n=1}^{\infty} |a_n(s) f(\sigma_n(s))|^2 \right)^{1/2}.$$

Since (R) implies that

$$\|\rho_s\| \leq |T|(|f|)(s),$$

we may apply the dominated convergence theorem, and it follows then that

$$\|f\|_1 \leq \left\| \left(\sum_{k=1}^{\infty} |a_k(s)|^2 |f(\sigma_k(s))|^2 \right)^{1/2} \right\|_1,$$

which is Lemma 1. □

Lemma II.2. Let $T : L^1([0, 1], m) \rightarrow L^1([0, 1], m)$ and $\alpha \geq 1$ be such that $\alpha \|f\|_1 \geq \|Tf\|_1 \geq \|f\|_1$ for all $f \in L^1$. Write

$$Tf(s) = \sum_{n=1}^{\infty} a_n(s) f(\sigma_n(s)) + \int_0^1 f(t) d\nu_s(t).$$

Then, for all $f \in L^1$,

$$(*) \quad \left\| \sum_{n=1}^{\infty} |a_n(s)| |f(\sigma_n(s))| \right\|_1 \leq \alpha \|f\|_1,$$

$$(**) \quad \alpha^{-1} \|f\|_1 \leq \left\| \max_n |a_n(s)| |f(\sigma_n(s))| \right\|_1.$$

Proof. It follows e.g. from (R) that $\| |T| \| = \|T\|$, hence $\| |T|(f) \|_1 \leq \alpha \|f\|_1$, which shows (*). For (**) we observe that, for any $s \in [0, 1]$,

$$\begin{aligned} & \left(\sum_{k=1}^{\infty} |a_k(s)|^2 |f(\sigma_k(s))|^2 \right)^{1/2} \\ & \leq \left(\sum_{k=1}^{\infty} |a_k(s)| |f(\sigma_k(s))| \right)^{1/2} \left(\max_k |a_k(s)| |f(\sigma_k(s))| \right)^{1/2}, \end{aligned}$$

hence by the Cauchy-Schwarz inequality

$$\begin{aligned} & \left\| \left(\sum_{k=1}^{\infty} |a_k(s)|^2 |f(\sigma_k(s))|^2 \right)^{1/2} \right\|_1 \\ & \leq \left\| \left(\sum_{k=1}^{\infty} |a_k(s)| |f(\sigma_k(s))| \right) \right\|_1^{1/2} \cdot \left\| \max_k |a_k(s)| |f(\sigma_k(s))| \right\|_1^{1/2}. \end{aligned}$$

It follows now from (*) and Lemma 1 that

$$\|f\|_1 \leq \sqrt{\alpha} \|f\|_1^{1/2} \left\| \max_k |a_k(s)| |f(\sigma_k(s))| \right\|_1^{1/2},$$

which shows (**). □

Lemma II.2 allows us to show the following theorem due to L. Dor ([7], Corollary 3.3) through a Hahn-Banach argument, distinct however from the original argument ([7]). Note that we do not have the better constant 1.6 obtained by L. Dor in the real case.

Theorem II.3. Let $T : L^1 \rightarrow L^1$ be an isomorphism into L^1 such that $\|T\| \cdot \|T^{-1}\| < \sqrt{2}$, with $T^{-1} : T(L^1) \rightarrow L^1$. Then $T(L^1)$ is a complemented subspace of L^1 .

Proof. We may and do assume that T satisfies the assumptions of Lemma II.2, with $\alpha < \sqrt{2}$. Define $\Psi : L^1 \rightarrow L^1(c_0)$ by

$$\Psi(f) = (|a_n(s)|f(\sigma_n(s)))_{n \geq 1}.$$

Note that, by $(*)$, Ψ actually takes its values in $L^1(\ell_1)$. By $(**)$, we have

$$\|\Psi(f)\|_{L^1(c_0)} \geq \alpha^{-1}\|f\|_1 \quad \text{for all } f \in L^1.$$

Hence, there exists $F \in L^1(c_0)^* = L^\infty(\ell_1)$, with $F \geq 0$, $\|F\| \leq \alpha$, such that

$$\int_0^1 f(s) ds = F[\Psi(f)] \quad \text{for all } f \in L^1.$$

That means there exists a sequence $(b_n)_{n \geq 1}$ of positive functions such that

$$\left\| \sum_{n=1}^{\infty} b_n(s) \right\|_{\infty} \leq \alpha,$$

and

$$\int_0^1 f(s) ds = \int_0^1 \sum_{n=1}^{\infty} b_n(s) |a_n(s)| f(\sigma_n(s)) ds \quad \text{for all } f \in L^1.$$

Choose disjoint measurable subsets E_n , $n \geq 1$, of $[0, 1] \times [0, 1]$ such that

$$m(\{x \mid (s, x) \in E_n\}) = \frac{b_n(s)}{\alpha} \quad \text{for all } s \in [0, 1], n \geq 1.$$

We define $U : L^1([0, 1]) \rightarrow L^1([0, 1]^2)$ by

$$Uf(s, x) = \sum_{n=1}^{\infty} a_n(s) f(\sigma_n(s)) \mathbf{1}_{E_n}(s, x).$$

For any $f \in L^1$, one has

$$\|Uf\|_1 = \int_0^1 \sum_{n=1}^{\infty} \frac{b_n(s)}{\alpha} |a_n(s)| |f(\sigma_n(s))| ds = \frac{1}{\alpha} \|f\|_1.$$

Hence $U(L^1)$ is isometric to L^1 , and it follows (see [22], Chap. 6, Section 17, Theorem 3) that $U(L^1)$ is 1-complemented in $L^1([0, 1]^2)$.

Define $\bar{T} : L^1([0, 1]) \rightarrow L^1([0, 1]^2)$ by

$$\bar{T}f(s, x) = Tf(s) = \sum_{n=1}^{\infty} a_n(s) f(\sigma_n(s)) + \int_0^1 f(u) d\nu_s(u).$$

If we denote

$$Uf(s, x) = \int_0^1 f(u) d\mu_{(s,x)}^U(u),$$

$$\bar{T}f(s, x) = \int_0^1 f(u) d\mu_{(s,x)}^{\bar{T}}(u),$$

then clearly

$$\mu_{(s,x)}^U = a_k(s) \delta_{\sigma_k(s)}$$

for the unique $k \geq 1$ such that $(s, x) \in E_k$, and

$$\mu_{(s,x)}^{\bar{T}} = \sum_{n=1}^{\infty} a_n(s) \delta_{\sigma_n(s)} + \nu_s,$$

for all (s, x) . Hence, we have

$$\mu_{(s,x)}^{\bar{T}-U} = \sum_{n \neq k} a_n(s) \delta_{\sigma_n(s)} + \nu_s,$$

and it easily follows that

$$|\bar{T}| - |U| = |\bar{T} - U| \geq 0.$$

Note that for all $f \in L^1$:

$$\| |U|(f) \|_1 = \|Uf\|_1 = \frac{1}{\alpha} \|f\|_1.$$

On the other hand,

$$\| |\bar{T}|(f) \|_1 = \| |T|(f) \|_1 \leq \alpha \|f\|_1.$$

Hence, for any $f \geq 0$,

$$\begin{aligned} \|\overline{T} - U\|(f) \|_1 &= \|\overline{T}(f) - |U|(f)\|_1 = \|\overline{T}(f)\|_1 - \||U|(f)\|_1 \\ &\leq (\alpha - \alpha^{-1})\|f\|_1 = (\alpha^2 - 1)\|Uf\|_1. \end{aligned}$$

Observe now that U maps disjoint functions to disjoint functions, hence if $f = f^+ - f^-$

$$\|Uf\|_1 = \|Uf^+\|_1 + \|Uf^-\|_1,$$

and thus for any $g \in L^1$:

$$(2) \quad \|\overline{T} - U\|g\|_1 \leq (\alpha^2 - 1)\|Ug\|_1 = \beta\|Ug\|_1.$$

Since $\alpha < \sqrt{2}$, we have $\beta = \alpha^2 - 1 < 1$. Let $\Pi : L^1([0, 1]^2) \rightarrow U(L^1)$ be a projection with $\|\Pi\| = 1$. It follows from (2) that $\Pi|_{\overline{T}(L^1)}$ is an isomorphism from $\overline{T}(L^1)$ onto $U(L^1)$.

Let $S = \Pi^{-1} : U(L^1) \rightarrow \overline{T}(L^1)$. It is easily seen that $S\Pi$ defines a projection from $L^1([0, 1]^2)$ onto $\overline{T}(L^1)$. Restricting $S\Pi$ to the functions which are independent of x , gives a projection from L^1 to $T(L^1)$. □

Our goal is now to provide a similar approach to Alspach’s result on near isometries ([1]). We first prove:

Theorem II.4. *Let $T : L^1 \rightarrow L^1$ be such that*

$$\alpha\|f\|_1 \geq \|Tf\|_1 \geq \|f\|_1 \quad \text{for all } f \in L^1.$$

Then for each $\varepsilon > 0$, there exists an operator $S_\varepsilon : L^1 \rightarrow L^1$ of the form

$$S_\varepsilon f(s) = \sum_{k=1}^N c_k(s) f(\sigma_k(s)),$$

and such that

$$\|T - S_\varepsilon\| \leq (\alpha - 1) + \varepsilon.$$

Proof. We consider again

$$\Psi(f) = (|a_n(s)| f(\sigma_n(s)))_{n \geq 1},$$

but this time as an operator from L^1 to $L^1(\ell_2)$. By Lemma 1, we have for all f

$$\|f\|_1 \leq \|\Psi(f)\|_{L^1(\ell_2)}.$$

Since $L^1(\ell_2)^* = L^\infty(\ell_2)$, there is $(b_n)_{n \geq 1}$ with $b_n \geq 0$ such that

$$\begin{cases} \left\| \left(\sum b_n^2 \right)^{1/2} \right\|_\infty \leq 1, \\ \int_0^1 \sum b_n(s) |a_n(s)| f(\sigma_n(s)) ds = \int_0^1 f ds, \end{cases}$$

for all f . We may and do redefine the a_n, b_n, σ_n 's in such a way that $b_1 \geq b_2 \geq \dots$; hence $0 \leq b_n(s) \leq 1/\sqrt{n}$ for all s and n . Then

$$\begin{aligned} & \int_0^1 \left| \sum_{n=N+1}^\infty b_n(s) a_n(s) f(\sigma_n(s)) \right| ds \\ & \leq \frac{1}{\sqrt{N}} \int_0^1 \sum_{n=N+1}^\infty |a_n(s)| |f(\sigma_n(s))| ds \leq \frac{1}{\sqrt{N}} \|T\| \cdot \|f\|_1, \end{aligned}$$

thus

$$\int_0^1 \sum_{n=1}^N b_n(s) |a_n(s)| |f(\sigma_n(s))| ds \geq \|f\|_1 - \frac{1}{\sqrt{N}} \|T\| \cdot \|f\|_1.$$

We let

$$S_N f(s) = \sum_{n=1}^N a_n(s) b_n(s) f(\sigma_n(s)).$$

For any s

$$\mu_s^{T-S_N} = \sum_{n=1}^N a_n(s) (1 - b_n(s)) \delta_{\sigma_n(s)} + \nu_s + \sum_{n=N+1}^\infty a_n(s) \delta_{\sigma_n(s)},$$

therefore $|T - S_N| = |T| - |S_N|$, and for any $f \geq 0$, one has

$$\| |T - S_N|(f) \|_1 = \| |T|(f) - |S_N|(f) \|_1 = \| |T|(f) \|_1 - \| |S_N|(f) \|_1.$$

Since

$$\| |T|(f) \|_1 \leq \alpha \|f\|_1 \quad \text{and} \quad \| |S_N|(f) \|_1 \geq \|f\|_1 - \frac{1}{\sqrt{N}} \|T\| \cdot \|f\|_1,$$

one has, for all $f \geq 0$,

$$\| |T - S_N|(f) \|_1 \leq (\alpha - 1) \|f\|_1 + \frac{1}{\sqrt{N}} \|T\| \cdot \|f\|_1,$$

and thus

$$\|T - S_N\| = \| |T - S_N| \| \leq (\alpha - 1) + \frac{\|T\|}{\sqrt{N}}.$$

The conclusion follows, since N is arbitrary. □

Proposition II.5. Let $S : L^1 \rightarrow L^1$ be an operator such that

$$Sf(s) = \sum_{k=1}^{\infty} a_k(s) f(\sigma_k(s)) \quad \text{and} \quad \alpha \|f\|_1 \geq \|Sf\|_1 \geq \|f\|_1$$

for every $f \in L^1$. Then there exists an operator $U : L^1 \rightarrow L^1$, of the form

$$Uf(s) = c(s)f(\sigma(s)),$$

such that

$$\|S - U\| \leq 2 \left(\alpha - \frac{1}{\alpha} \right).$$

Proof. We consider, as before, $\Psi(f) = (|a_n(s)| f(\sigma_n(s)))_{n \geq 1}$ as an operator into $L^1(c_0)$. Since by (**)

$$\alpha^{-1} \|\Psi(f)\|_{L^1(c_0)} \leq \|f\|_{L^1},$$

there exists $(b_k)_k$ with $b_k \geq 0$, $\sum_k b_k \leq 1$, and

$$\int_0^1 \sum_{k=1}^{\infty} |a_k(s)| b_k(s) f(\sigma_k(s)) ds = \frac{1}{\alpha} \int_0^1 f(s) ds.$$

We may and do assume that $b_1 \geq b_2 \geq \dots$. Therefore, for all $k \geq 2$,

$$1 - b_k \geq 1 - b_1 \geq b_k.$$

Since $\alpha \|f\|_1 \geq \| |S|(f) \|_1$, one has, for all $f \geq 0$,

$$\begin{aligned} \left(\alpha - \frac{1}{\alpha}\right) \|f\|_1 &\geq \int_0^1 \sum_{k=1}^{\infty} |a_k(s)| (1 - b_k(s)) f(\sigma_k(s)) ds \\ &\geq \int_0^1 \sum_{k=2}^{\infty} |a_k(s)| b_k(s) f(\sigma_k(s)) ds. \end{aligned}$$

Hence if we define

$$Vf(s) = \sum_{k=2}^{\infty} a_k(s) b_k(s) f(\sigma_k(s))$$

and

$$\begin{aligned} Wf(s) &= \sum_{k=1}^{\infty} a_k(s) b_k(s) f(\sigma_k(s)) \\ &= a_1(s) b_1(s) f(\sigma_1(s)) + Vf(s), \end{aligned}$$

one has, for all $f \geq 0$,

$$\| |S - W|(f) \|_1 = \| |S|(f) - |W|(f) \|_1 \leq \left(\alpha - \frac{1}{\alpha}\right) \|f\|_1,$$

and also

$$\| |V|(f) \|_1 \leq \left(\alpha - \frac{1}{\alpha}\right) \|f\|_1,$$

thus

$$\|S - W\| \leq \left(\alpha - \frac{1}{\alpha}\right) \quad \text{and} \quad \|V\| \leq \left(\alpha - \frac{1}{\alpha}\right).$$

We let now $U = W - V$. We have:

$$\|S - U\| = \|S - W + V\| \leq \|S - W\| + \|V\| \leq 2 \left(\alpha - \frac{1}{\alpha}\right). \quad \square$$

Proposition II.6. Let $U : L^1 \rightarrow L^1$ be of the form

$$Uf(s) = c(s)f(\sigma(s)),$$

and assume that for some $\alpha \geq 1$ we have, for all $f \in L^1$,

$$\alpha \|f\|_1 \geq \|U(f)\|_1 \geq \|f\|_1.$$

Then there exists an isometry J , from L^1 to L^1 , such that $\|U - J\| \leq (\alpha - 1)$.

Proof. By the Radon-Nikodym theorem, there exists w such that, for all $f \in L^1$,

$$\int_0^1 |U|(f)(s) ds = \int_0^1 w(s)f(s) ds.$$

For all $f \geq 0$, one has

$$\alpha \int_0^1 f(s) ds \geq \int_0^1 w(s)f(s) ds \geq \int_0^1 f(s) ds.$$

It follows that $1 \leq w(s) \leq \alpha$ for all s . We define

$$Jf(s) = \frac{c(s)}{w(\sigma(s))} f(\sigma(s)),$$

and compute

$$\begin{aligned} \|Jf\|_1 &= \int_0^1 |Jf(s)| ds = \int_0^1 \frac{|c(s)|}{w(\sigma(s))} |f(\sigma(s))| ds \\ &= \int_0^1 |c(s)| \left(\frac{|f|}{w} \right) (\sigma(s)) ds \\ &= \int_0^1 |U| \left(\frac{|f|}{w} \right) (s) ds = \int_0^1 w(s) \left(\frac{|f|}{w} \right) (s) ds = \|f\|_1. \end{aligned}$$

Hence J is an isometry. Moreover, for any s , we have

$$\mu_s^{|U-J|} = |c(s)| \left(1 - \frac{1}{w(\sigma(s))} \right) \delta_{\sigma(s)},$$

and since

$$0 \leq 1 - \frac{1}{w(s)} \leq 1 - \frac{1}{\alpha} = \frac{\alpha - 1}{\alpha},$$

we have

$$\|U - J\| = \| |U - J| \| \leq \frac{\alpha - 1}{\alpha} \|U\| \leq (\alpha - 1). \quad \square$$

From Theorem II.4, Proposition II.5, and Proposition II.6 we deduce the following quantitative improvement of Alspach's theorem ([1]).

Theorem II.7. *There is a function $\varphi(\alpha)$, with $\lim_{\alpha \rightarrow 1^+} \varphi(\alpha) = 0$ and such that, if $T : L^1 \rightarrow L^1$ satisfies for some $\alpha \geq 1$*

$$\alpha \|f\|_1 \geq \|Tf\|_1 \geq \|f\|_1 \quad \text{for all } f \in L^1,$$

then there is an isometry $J : L^1 \rightarrow L^1$ such that, for all $f \in L^1$,

$$\|T - J\| \leq \varphi(\alpha).$$

Moreover, for $1 \leq \alpha \leq 1.08$ we can choose $\varphi(\alpha) = 13(\alpha - 1)$.

Proof. The first part clearly follows from Theorem II.4, Proposition II.5, and Proposition II.6. For the second part, apply Theorem II.4 to find an S so that $\|T - S\| \leq \frac{17}{16}(\alpha - 1) = \delta_1 < 1$. Then $S' = S/(1 - \delta_1)$ verifies $\|f\|_1 \leq \|S'f\|_1 \leq \alpha_1 \|f\|_1$, with

$$\alpha_1 = \frac{\delta_1 + \alpha}{1 - \delta_1} \leq 1.28 \quad \text{and} \quad \|T - S'\| \leq \frac{\delta_1}{1 - \delta_1}(1 + \alpha) \leq 2.5(\alpha - 1).$$

Now, Proposition II.5 gives an operator U so that $\|S' - U\| \leq 2(\alpha_1 - \alpha_1^{-1}) = \delta_2 < 1$; the operator $U' = U/(1 - \delta_2)$ verifies $\|f\|_1 \leq \|U'f\|_1 \leq \alpha_2 \|f\|_1$, with $\alpha_2 = (\delta_2 + \alpha_1)/(1 - \delta_2)$, and we have

$$\|S' - U'\| \leq \frac{\delta_2}{1 - \delta_2}(1 + \alpha_1) \leq 3.4(\alpha - 1).$$

Finally, Proposition II.6 gives a J so that $\|U' - J\| \leq \alpha_2 - 1 \leq 6.8(\alpha - 1)$. \square

Remark II.8. We give a proof of Theorem II.7 in the special case where $T = C_\mu$ is the convolution by a measure μ on $L^1(\mathbb{T})$. Suppose that $\|\mu\| = 1$ and that $\|f * \mu\|_1 \geq (1 - \varepsilon)\|f\|_1$. Write the atomic part of μ as $\sum_{k=1}^{+\infty} a_k \delta_{x_k}$. We have

$$(3) \quad \sum_{k=1}^{+\infty} |a_k| \leq \|\mu\| = 1.$$

On the other hand, Wiener's theorem ([20], p. 42; [13], p. 415):

$$\lim_{N \rightarrow +\infty} \frac{1}{2N + 1} \sum_{n=-N}^N |\hat{\mu}(n)|^2 = \sum_{k=1}^{+\infty} |a_k|^2$$

says that

$$(4) \quad (1 - \varepsilon) \leq \sum_{k=1}^{+\infty} |a_k|^2,$$

since

$$|\hat{\mu}(n)| = \|\mu * e_n\|_1 \geq (1 - \varepsilon),$$

where $e_n(t) = e^{2\pi int}$. It follows from (3) and (4) that there exists an index K such that $|a_K| \geq (1 - 2\varepsilon)$; thus the distance between C_μ and translation by x_K is less than 4ε , and this concludes the proof.

There is an interesting link between Wiener’s theorem and the equation (1) in the proof of Lemma II.1. Indeed, let μ be a measure on the Cantor group $G = \{-1, 1\}^{\mathbb{N}}$, and let (W_k) and (D_k) be the Walsh functions and the elementary open sets with index $k \in \{-1, 1\}^{[\mathbb{N}]}$. With this notation one has, by Parseval’s formula,

$$2^{-l} \sum_{k \in 2^l} \left(\int W_k d\mu \right)^2 = \sum_{k \in 2^l} \mu(D_k)^2$$

for any $l \geq 1$, and thus equation (1) appears, through the natural measure preserving isomorphism between G and the unit interval, as the exact analogue for the Cantor group of Wiener’s theorem. It would be nice to state a generalization of Wiener’s theorem to arbitrary compact abelian groups.

III. Small subspaces of L^1 . In this section, we introduce and study the notion of *small subspace* of L^1 , which will be essential in the next section. In the sequel the notation $A \subsetneq D$ will be used to say that there is a real number $k > 0$ such that $kA \subseteq D$. If X is a subspace of L^1 , we denote by $C_X = \overline{B_X}^{\tau_m}$ the closure of its unit ball for the topology τ_m of convergence in measure. Recall that Bukhvalov and Lozanovski ([4]) showed that $P(B_{X^{\perp\perp}}) = C_X$, where P is the natural projection from L^{1**} to L^1 . We also denote, for $A \subseteq \Omega$,

$$L^1(A) = \{f \in L^1 \mid \mathbf{1}_{\Omega \setminus A} \cdot f = 0 \text{ a.e.}\}.$$

Definition III.1. A subspace X of $L^1(\Omega, m)$ is said to be *small* if there is no $A \subseteq \Omega$ of positive measure such that $B_{L^1(A)} \subsetneq \mathbf{1}_A \cdot B_X$.

In other words, the projection $f \mapsto f \cdot \mathbf{1}_A$ never maps X onto $L^1(A)$. It is easy to see (by lifting finite trees for instance) that reflexive subspaces of

L^1 are small. However, smallness is not an intrinsic notion, but is related to the position of X in L^1 . Indeed, let $Q : \ell_1 \rightarrow L^1$ be a quotient map, and $J : \ell_1 \rightarrow L^1$ be an isometric embedding. Then the range of $(J, Q) : \ell_1 \rightarrow L^1 \times L^1$ is isomorphic to ℓ_1 , but is not small. On the other hand, in $L^1([0, 1] \times [0, 1])$, the subspace of the functions which do not depend of the second coordinate is small, although it is isometric to L^1 . We are now going to give some examples and properties of small subspaces.

It will be useful to know that small subspaces actually satisfy a formally stronger property.

Proposition III.2. *If X is a small subspace of L^1 , there is no $E \subseteq \Omega$ of positive measure such that $B_{L^1(E)} \subsetneq \mathbf{1}_E \cdot C_X$.*

Proof. Assume that $B_{L^1(E)} \subseteq \mathbf{1}_E \cdot (kC_X)$, and let $\{x_n\}$ be dense in $B_{L^1(E)}$. There are $f_n \in kB_X$, and $E_n \subseteq E$, with $m(E_n) \leq 4^{-n}m(E)$, such that $|x_n - f_n| \leq 1/n$ on $(E \setminus E_n)$. If now $A = E \setminus (\bigcup_{n \geq 1} E_n)$, we have $B_{L^1(A)} \subseteq \mathbf{1}_A \cdot (kB_X)$, and so X is not small. \square

The following notion is closely related to smallness. It should be however noted that, except in the proof of Proposition III.9, only the property stated in Proposition III.4 will actually be used.

Definition III.3. Let T be an operator on L^1 , written as

$$Tf(x) = \int_0^1 f(s) d\mu_x(s).$$

T will be called a *strong Enflo operator* if there exists a set of positive measure on which $\mu_s(\{s\}) \neq 0$.

In other words, T has a “diagonal part”. Such operators necessarily have an atomic part, and so are isomorphisms on some subspace $L^1(A)$ ([18], Theorem 5.5); this means that they are Enflo operators ([8], Theorem 4.1). Translations are Enflo operators but not strong. The proof of Theorem 5.4 in [18] essentially shows:

Proposition III.4. *If T is a strong Enflo operator, there exists a set A of positive measure such that*

$$B_{L^1(A)} \subsetneq \mathbf{1}_A \cdot T(B_{L^1(A)}).$$

Indeed, if T is a strong Enflo operator and $T^a f(s) = \sum_{n=1}^{+\infty} a_n(s) f(\sigma_n(s))$ represents the atomic part of T , there exist an n and a set C of positive measure on which $a_n(s) \neq 0$ and $\sigma_n(s) = s$. It follows that, with the notation used in the proof of Theorem 5.4 of [18], $\sigma_n(B_{m,i} \cap B) \subseteq C_{m,i}$, and this set $A = B_{m,i} \cap B$ gives the conclusion.

Lemma III.5. *Let X be a subspace of $L^1(\Omega)$, and $T : L^1(A) \rightarrow L^1(\Omega)$ an operator, where $A \subseteq \Omega$. Then $T(B_{L^1(A)}) \subseteq \mathbf{1}_A \cdot C_X$ if and only if there exists an operator $\tilde{T} : L^1(A) \rightarrow X^{\perp\perp}$ such that $T = \mathbf{1}_A P \tilde{T}$.*

Proof. The sufficiency comes directly from Bukhvalov-Lozanovski’s theorem $P(B_{X^{\perp\perp}}) = C_X$. To see the other direction, write $L^1(A) = \bigcup_{n \geq 1} L_n$ with $L_n \subseteq L_{n+1}$ and L_n isometric to $\ell_1^{\dim L_n}$. Since $T(B_{L^1(A)}) \subseteq \mathbf{1}_A \cdot C_X$, there exist operators $T_n : L_n \rightarrow X$, such that $\|T_n\| \leq M$ and $d_m(\mathbf{1}_A T_n f, T f) \leq 2^{-n}$ for $f \in B_{L_n}$, where d_m is a distance generating the topology τ_m . Taking a w^* -limit of the T_n ’s along an ultrafilter, gives $\tilde{T} : L^1(A) \rightarrow X^{\perp\perp}$ such that $\mathbf{1}_A P \tilde{T} = T$. Indeed, if $F \in L^{1**}$ is the w^* -limit of a bounded filter $(f_\alpha)_\alpha$ of functions in L^1 , there exists a sequence $(c_k)_k$ of convex combinations of the f_α ’s, such that $d_m(c_k, P(F)) \xrightarrow{k \rightarrow +\infty} 0$ (see [16], Lemma IV.3.1). \square

We can now state the following characterization of small subspaces.

Proposition III.6. *A subspace X of L^1 is small if and only if no strong Enflo operator T verifies $T(B_{L^1}) \subseteq C_X$.*

Proof. Suppose that T is a strong Enflo operator such that $T(B_{L^1}) \subseteq C_X$. Since, by Proposition III.4, we have $B_{L^1(A)} \subseteq \mathbf{1}_A \cdot T(B_{L^1(A)})$ for some set A of positive measure, we obtain $B_{L^1(A)} \subseteq \mathbf{1}_A \cdot C_X$, in contradiction with the smallness of X . Suppose now that X is not small, and let A be of positive measure and such that $B_{L^1(A)} \subseteq \mathbf{1}_A \cdot C_X$. By applying Lemma III.5 to the natural injection $j_A : L^1(A) \rightarrow L^1$, we obtain $\tilde{j}_A : L^1(A) \rightarrow X^{\perp\perp}$, such that $j_A = \mathbf{1}_A P \tilde{j}_A$. Defining $T f = P \tilde{j}_A(\mathbf{1}_A f)$ for $f \in L^1$, we get a strong Enflo operator $T : L^1 \rightarrow L^1$, since, by uniqueness of the representation of T , we have $\mu_s^T = \delta_s$ for almost every $s \in A$, and we have also $T(B_{L^1(A)}) \subseteq C_X$. \square

Subspaces X for which $C_X = B_X$, i.e., nicely placed subspaces (see [10], Definition 2), are of particular interest: if X is nicely placed and no operator from L^1 to X is strongly Enflo, then X is small. Since, by Proposition III.4,

the range of every strong Enflo operator contains an isomorphic copy of L^1 , we get:

Proposition III.7. Every nicely placed subspace of L^1 which contains no subspace isomorphic to L^1 is small.

Note that the above condition is not necessary. For instance, the subspace of $L^1([0, 1]^2)$ consisting of the functions which do not depend on the second variable is isometric to L^1 and nicely placed, but it is small. A special case of Proposition III.7 is:

Proposition III.8. For every nicely placed Riesz subset Λ of \mathbb{Z} , $L^1_\Lambda(\mathbb{T})$ is small.

Recall that Λ is a Riesz set if every measure whose Fourier transform is carried by Λ is absolutely continuous, and it is nicely placed if L^1_Λ is nicely placed in L^1 (see [11], Definition 1.4 or [16], Definition IV.4.2). Proposition III.8 follows from the fact that Λ is a Riesz set if and only if $L^1_\Lambda(\mathbb{T})$ has the Radon-Nikodym Property ([27]). A stronger result will be shown below (Proposition III.10).

Proposition III.6 gives also a stability result.

Proposition III.9. Let X and Y be two small subspaces such that $X \cap Y = \{0\}$ and $X \oplus Y$ is closed. Then $X \oplus Y$ is small.

Proof. Let $T : L^1 \rightarrow L^1$ be such that $T(B_{L^1}) \subseteq C_{X \oplus Y}$. Lemma III.5 with $A = \Omega$ gives $\tilde{T} : L^1 \rightarrow (X \oplus Y)^{\perp\perp} = X^{\perp\perp} \oplus Y^{\perp\perp}$ such that $T = P\tilde{T}$. Let $\pi : X^{\perp\perp} \oplus Y^{\perp\perp} \rightarrow X^{\perp\perp}$ be the projection. Writing $U = \pi\tilde{T}$ and $V = (1 - \pi)\tilde{T}$, we have $T = PU + PV$, and $PU(B_{L^1}) \subseteq C_X$ and $PV(B_{L^1}) \subseteq C_Y$. By Proposition III.6, PU and PV cannot be strong Enflo operators, so $\mu_s^{PU}(\{s\}) = \mu_s^{PV}(\{s\}) = 0$ for almost all s . Hence $\mu_s^T(\{s\}) = 0$ for almost all s , and T is not a strong Enflo operator. By Proposition III.6 again, we can conclude that $X \oplus Y$ is small. \square

We now examine the translation invariant case. A subset Λ of \mathbb{Z} is said to be *semi-Riesz* ([37]) if every measure σ whose Fourier transform is carried by Λ has no atomic part. Every Riesz set is semi-Riesz, and Wiener's theorem ([20], p. 42; [13], p. 415):

$$\lim_{N \rightarrow +\infty} \frac{1}{2N + 1} \sum_{n=-N}^N |\hat{\sigma}(n)|^2 = \sum_{t \in \mathbb{T}} |\sigma(\{t\})|^2$$

shows that every set Λ with density zero, that is

$$d(\Lambda) = \lim_{N \rightarrow +\infty} \frac{\#(\Lambda \cap \{-N, \dots, N\})}{2N + 1} = 0,$$

is semi-Riesz. The set $\Lambda = \{\sum_{k=1}^n \varepsilon_k 4^k \mid \varepsilon_k = -1, 0, 1, n \geq 1\}$ is therefore semi-Riesz and not Riesz ([37], p. 126). We have:

Proposition III.10. *Let Λ be a nicely placed subset of \mathbb{Z} . Then $L^1_\Lambda(\mathbb{T})$ is small if and only if Λ is a semi-Riesz set.*

Proof. Suppose that $L^1_\Lambda(\mathbb{T})$ is not small, and let $T : L^1(\mathbb{T}) \rightarrow L^1(\mathbb{T})$ be the operator constructed in the proof of Proposition III.6. We have readily

$$\begin{cases} \mu_s^T(\{s\}) = 0 & \text{if } s \notin A, \\ \mu_s^T(\{s\}) = 1 & \text{if } s \in A. \end{cases}$$

Let $\tau_t : L^1(\mathbb{T}) \rightarrow L^1(\mathbb{T})$ be the translation given by

$$\tau_t f(u) = f(t + u).$$

Since for every operator U we have

$$\mu_s^{\tau_t^{-1} U \tau_t}(\{s\}) = \mu_{s-t}^U(\{s-t\}),$$

if we define $\tilde{T} : L^1(\mathbb{T}) \rightarrow L^1(\mathbb{T})$ by

$$\tilde{T} = \int_{\mathbb{T}} \tau_t^{-1} T \tau_t \, dm(t),$$

we have, for all s ,

$$\mu_s^{\tilde{T}}(\{s\}) = \int_{\mathbb{T}} \mu_{s-t}^T(\{s-t\}) \, dm(t),$$

and so $\mu_s^{\tilde{T}}(\{s\}) = m(A)^{-1}$. Since $\tau_t \tilde{T} = \tilde{T} \tau_t$ for every $t \in \mathbb{T}$, there exists a measure σ such that $\tilde{T}(f) = \sigma * f$. We have then $\sigma(\{0\}) = m(A)^{-1}$, so the atomic part of σ is non zero, although $\sigma \in \mathcal{M}_\Lambda$. Conversely, suppose that

L^1_Λ is small. Let σ be a measure with spectrum in Λ , and with atomic part

$$\sigma_a = \sum_{n=1}^{+\infty} a_n \delta_{x_n}.$$

The convolution operator defined by $Tf = \sigma * f$ maps $L^1(\mathbb{T})$ into L^1_Λ , so that $T(B_{L^1(\mathbb{T})}) \subseteq B_{L^1_\Lambda}$. On the other hand, its representing measure has the atomic part

$$\mu_s^{T_a} = \sum_{n=1}^{+\infty} a_n \delta_{s-x_n}.$$

Since L^1_Λ is small and nicely placed, T cannot be a strong Enflo operator, so $\mu_s^{T_a}(\{s\}) = 0$ almost everywhere. This is possible only if $x_n \neq 0$ or $a_n = 0$. Considering now, for each $p \geq 1$, the translated measure $\sigma * \delta_{-x_p}$ (which still have its spectrum in Λ) instead of σ , gives that $a_p = 0$, since its atomic part is $\sum_{n=1}^{+\infty} a_n \delta_{x_n-x_p}$. Hence $\sigma_a = 0$, and we are done. \square

Example. Y. Meyer ([28], Theorem 6) showed that, if $(n_k)_{k \geq 1}$ grows fastly enough:

$$n_{k+1} > 6(n_1 + \dots + n_k) \quad \text{and} \quad \sum_{k=1}^{+\infty} \frac{n_k}{n_{k+1}} < +\infty,$$

then the set

$$\Lambda = \left\{ \sum_{k=1}^K \varepsilon_k n_k \mid \varepsilon_k = -1, 0, 1, K \geq 1 \right\}$$

is closed for the Bohr topology; hence it is nicely placed. Since it has density zero, $L^1_\Lambda(\mathbb{T})$ is small, but Λ is not a Riesz set, as the spectrum of a Riesz product. F. Parreau pointed out to us that the same conclusion holds for

$$\Lambda = \left\{ \sum_{k=1}^K \varepsilon_k 4^k \mid \varepsilon_k = -1, 0, 1, K \geq 1 \right\}.$$

We do not know whether L^1_Λ contains a subspace isomorphic to L^1 whenever Λ is the spectrum of a Riesz product.

For every subset Λ of \mathbb{Z} , denote by $[\Lambda]$ the smallest nicely placed subset of \mathbb{Z} containing Λ (see [11], p. 306). $[\Lambda]$ is contained in the Bohr-closure of Λ ([11], Corollary 2.6).

Corollary III.11. If $L^1_\Lambda(\mathbb{T})$ is not small, $[\Lambda]$ contains a translate of the spectrum of a Riesz product.

Proof. If $L^1_\Lambda(\mathbb{T})$ is not small, $L^1_{[\Lambda]}(\mathbb{T})$ is a fortiori not small. Hence there is a measure σ with spectrum in $[\Lambda]$ which has an atomic part, and Wiener’s theorem shows that $\hat{\sigma}(n) \not\rightarrow_{|n| \rightarrow +\infty} 0$. Host-Parreau’s theorem ([17]) ends the proof. □

IV. Lifting of operators between quotients by small subspaces of L^1 . In this section, we show that operators between quotients of L^1 by small subspaces lift to operators on L^1 . We begin with a special case, which is much easier than the general one. The following notion will be used:

Definition IV.1. The operator $T : L^1 \rightarrow L^1$ is said to be a *Daugavet operator* if the Daugavet equation $\|I + \lambda T\| = 1 + \|T\|$ is fulfilled for every λ with $|\lambda| = 1$.

A. Plichko and M. Popov showed ([30], Section 9, Theorem 8) that every non-Enflo operator is Daugavet. In fact, more is true: T is a Daugavet operator whenever it is not strongly Enflo. Indeed, denoting the representing measures of T by $\mu_s, s \in [0, 1]$, we have

$$\|T\| = \sup_E \frac{1}{m(E)} \int_0^1 |\mu_s|(E) dm(s),$$

and

$$\|I + \lambda T\| = \sup_E \frac{1}{m(E)} \int_0^1 |\delta_s + \lambda \mu_s|(E) dm(s).$$

Since T is not strongly Enflo, we have $\mu_s(\{s\}) = 0$ for almost all s , and the measures δ_s and μ_s are disjoint for these values. Hence, if for every $\varepsilon > 0$ we choose a measurable set E such that

$$\|T\| \leq \frac{1}{m(E)} \int_0^1 |\mu_s|(E) dm(s) + \varepsilon,$$

we get, since we may assume that $s \in E$,

$$1 + \|T\| \leq \frac{1}{m(E)} \int_0^1 |\delta_s + \lambda \mu_s|(E) dm(s) + \varepsilon \leq \|I + \lambda T\| + \varepsilon,$$

which gives the result. D. Werner used a similar argument in ([36]).

Proposition IV.2. Let X and Y be nicely placed subspaces of L^1 . Assume that X and Y are small, and that

$$\text{dist}(L^1/X, L^1/Y) < 2.$$

Then there exists an invertible operator $V : L^1 \rightarrow L^1$ such that $V(X) = Y$. Moreover, for $0 \leq \alpha < 1$, if

$$\text{dist}(L^1/X, L^1/Y) = 1 + \alpha,$$

there is such a V satisfying $\|V\| \cdot \|V^{-1}\| \leq (1 + \alpha)/(1 - \alpha)$.

Proof. Let S be an isomorphism from L^1/X onto L^1/Y such that $\|S\| \cdot \|S^{-1}\| < 2$. By the lifting property of nicely placed subspaces (see [23], Theorem 1, [16], Proposition IV.2.12, or [19], Proposition 2.1), there exist U and V such that

$$\begin{cases} Q_X U = S^{-1} Q_Y, \\ Q_Y V = S Q_X, \end{cases}$$

and such that $\|U\| = \|S^{-1}\|$ and $\|V\| = \|S\|$. Hence the following diagram is commutative:

$$\begin{array}{ccc} L^1 & \begin{array}{c} \xrightarrow{V} \\ \xleftarrow{U} \end{array} & L^1 \\ \begin{array}{c} \downarrow Q_X \\ \\ \downarrow \end{array} & & \begin{array}{c} \downarrow Q_Y \\ \\ \downarrow \end{array} \\ L^1/X & \begin{array}{c} \xrightarrow{S} \\ \xleftarrow{S^{-1}} \end{array} & L^1/Y \end{array}$$

For every $f \in L^1$, we have

$$Q_Y V U f = S Q_X U f = S S^{-1} Q_Y f = Q_Y f,$$

hence

$$(I - VU)(L^1) \subseteq Y,$$

and similarly $(I - UV)(L^1) \subseteq X$. Since X and Y are nicely placed and small, the operators $I - UV$ and $I - VU$ are not strongly Enflo, and so are Daugavet operators. Therefore, we have

$$\|I - UV\| = \|UV\| - 1 \quad \text{and} \quad \|I - VU\| = \|VU\| - 1.$$

Since $\|UV\| \leq \|U\| \cdot \|V\| = \|S\| \cdot \|S^{-1}\| < 2$, it follows that U and V are invertible. Consequently $V(X) = Y$ since, for $f, g \in L^1$ such that $V(f) = g$, we have

$$\begin{aligned} g \in Y &\iff Q_Y(g) = 0 \iff Q_YV(f) = 0 \\ &\iff SQ_X(f) = 0 \iff Q_X(f) = 0 \iff f \in X. \end{aligned}$$

To conclude the proof, we observe that $\|V\| = \|S\|$ and

$$V^{-1} = \left[\sum_{k=0}^{+\infty} (I - UV)^k \right] U.$$

Hence

$$\|V^{-1}\| \leq \frac{\|S^{-1}\|}{2 - \|S\| \cdot \|S^{-1}\|},$$

since $\|U\| = \|S^{-1}\|$ and $\|I - UV\| \leq \|S\| \cdot \|S^{-1}\|$, and so V works. □

Remark. Note that Proposition IV.2 applies in particular to nicely placed subspaces which do not contain L^1 (by Proposition III.7). If L^1/X and L^1/Y are isometric, we have $\alpha = 0$ and, by the above, V is actually an invertible isometry of L^1 .

We now give the general statement.

Theorem IV.3. *Let X and Y be small subspaces of L^1 , and suppose that there is an isomorphism $S : L^1/X \rightarrow L^1/Y$ such that $\|S\|$ and $\|S^{-1}\|$ are less than $1 + \delta$, with $\delta < \frac{1}{25}$. Then, there exists an invertible operator $U : L^1 \rightarrow L^1$ such that $\|U\| \cdot \|U^{-1}\| \leq (1 + \delta)/(1 - 25\delta)$, and $d_{\mathcal{H}}(U(B_X), B_Y) \leq 71\delta/(1 - 25\delta)$, where $d_{\mathcal{H}}$ denotes the Hausdorff distance.*

Proof. Denote by Q_X and Q_Y the projections onto the respective quotient spaces. We first construct a bidual diagram

$$\begin{array}{ccc} L^{1**} & \begin{array}{c} \xrightarrow{W_1} \\ \xleftarrow{W_2} \end{array} & L^{1**} \\ \begin{array}{c} \downarrow Q_X^{**} \\ \\ \downarrow \end{array} & & \begin{array}{c} \downarrow Q_Y^{**} \\ \\ \downarrow \end{array} \\ L^{1**}/X^{\perp\perp} & \begin{array}{c} \xrightarrow{S^{**}} \\ \xleftarrow{S^{-1**}} \end{array} & L^{1**}/Y^{\perp\perp} \end{array}$$

Write $L^1 = \overline{\bigcup_{n \geq 1} L_n}$, where (L_n) is an increasing sequence of finite-dimensional spaces isometric to $\ell_1^{\dim(L_n)}$. For all $n \geq 1$, there exists $T_n : L_n \rightarrow L^1$ such that

$$Q_Y T_n = S Q_X|_{L_n},$$

and $\|T_n\| \leq (1 + 1/n) \|S\|$.

Lemma IV.4. Let \mathcal{U} be a free ultrafilter on the integers. There exists a (w^*-w^*) -continuous operator $W_1 : L^{1**} \rightarrow L^{1**}$ such that

$$\begin{cases} \|W_1\| \leq \|S\|, \\ Q_Y^{**} W_1 = S^{**} Q_X^{**}, \\ W_1(f) = F, \end{cases}$$

for any $f \in \bigcup_{n \geq 1} L_n$ and $F = w^*\text{-}\lim_{\mathcal{U}} T_n(f)$.

Proof. Define $w_1 : \bigcup_n L_n \rightarrow L^{1**}$ by

$$w_1(f) = w^*\text{-}\lim_{\mathcal{U}} T_n(f).$$

Clearly $\|w_1\| \leq \|S\|$, and w_1 can be extended to an operator $\bar{w}_1 : L^1 \rightarrow L^{1**}$. If now $\Pi : (L^1)^{(4)} \rightarrow L^{1**}$ denotes the canonical (w^*-w^*) -continuous projection (“restriction to L^{1**} ”), then $W_1 = \Pi \bar{w}_1^{**}$ works. Indeed, for any $g \in L^1$, one has

$$Q_Y^{**} \bar{w}_1(g) = S Q_X(g),$$

and the equation $Q_Y^{**} W_1 = S^{**} Q_X^{**}$ follows by w^* -continuity of all relevant operators. □

We get of course, by the same token, an operator W_2 with $Q_X^{**} W_2 = S^{-1**} Q_Y^{**}$ and similar properties.

Lemma IV.5. Let $P : L^{1**} \rightarrow L^1$ be the natural projection. Then

$$[P(I - W_1 W_2)](B_{L^{1**}}) \subseteq \|I - W_1 W_2\| \overline{B_Y}^{\tau^m}.$$

Proof. For any $h \in L^{1**}$, we have

$$Q_Y^{**} W_1 W_2 h = S^{**} Q_X^{**} W_2 h = S^{**} S^{-1**} Q_Y^{**} h = Q_Y^{**} h.$$

Hence $Q_Y^{**}(I - W_1W_2) = 0$ and

$$(I - W_1W_2)(B_{L^{**}}) \subseteq \|I - W_1W_2\|B_{Y^{\perp\perp}}.$$

The result follows by Bukhvalov-Lozanovski's theorem $P(B_{Y^{\perp\perp}}) = \overline{B_Y}^{\tau_m}$. \square

Then, writing $T = P(I - W_1W_2)|_{L^1}$, it follows from Lemma IV.5, the smallness of Y , and Proposition III.6 that T is not a strong Enflo operator; hence $-T$ satisfies the Daugavet equation, and we have

$$\begin{aligned} 1 + \|P(I - W_1W_2)|_{L^1}\| &= 1 + \|T\| = \|I - T\| = \|(PW_1W_2)|_{L^1}\| \\ &\leq \|S\| \cdot \|S^{-1}\| \leq 1 + \delta. \end{aligned}$$

Therefore

$$(1) \quad \|P(I - W_1W_2)|_{L^1}\| \leq \delta,$$

and hence, for any $f \in L^1$, one has

$$\|P(I - W_1W_2)(f)\|_1 \geq (1 - \delta)\|f\|_1.$$

On the other hand, one also has

$$\|W_1W_2f\|_1 \leq \|S\| \cdot \|S^{-1}\| \|f\|_1 \leq (1 + \delta)\|f\|_1.$$

Since P is an L -projection, we have then

$$\|W_1W_2f - PW_1W_2f\|_1 = \|W_1W_2f\|_1 - \|PW_1W_2f\|_1 \leq 2\delta\|f\|_1.$$

It follows from this and (1) that

$$\|(I - W_1W_2)|_{L^1}\| \leq 3\delta.$$

Since W_1 and W_2 are w^* -continuous, this gives

$$(2) \quad \|I - W_1W_2\| \leq 3\delta.$$

An identical argument shows that

$$(3) \quad \|I - W_2W_1\| \leq 3\delta.$$

Since $\delta < \frac{1}{3}$, (2) and (3) show that W_1W_2 and W_2W_1 are invertible. Hence W_1 and W_2 are themselves invertible.

Lemma IV.6. In the notation of Lemma IV.4, one has

$$\| \text{-diam}(L_f) \leq 6\delta \|W_2^{-1}\| \|f\|_1,$$

where L_f is the set of all w^* -cluster points of the $T_n f$'s.

Proof. Since $\|f - W_2W_1f\|_1 \leq 3\delta \|f\|_1$ by (2), one has

$$\|W_2^{-1}f - W_1f\|_1 \leq 3\delta \|W_2^{-1}\| \|f\|_1.$$

By Lemma IV.4, any $F \in L_f$ can be written $F = W_1(f)$ for an appropriate W_1 , and the conclusion follows. \square

Note that it follows, from (2) and (3), that

$$\|W_2^{-1}\| \leq \psi(\delta) = \frac{1 + \delta}{1 - 3\delta}.$$

Lemma IV.7. Let Z be a Banach space such that $Z^{**} = Z \oplus_1 Z_s$. Let $(z_n)_n$ be a sequence in Z and L be the set of all its w^* -cluster points. If

$$\| \text{-diam}(L) = \varepsilon,$$

we have $\text{dist}(l, Z) \leq \varepsilon$ for all $l \in L$.

Proof. Let $l \in L$. By translating if necessary, we may and do assume that $l \in Z_s$. Let $K = (B_{Z^*}, w^*)$. By ([6], Lemma III.2.4), we have

$$\|l\| = \text{dist}(l, Z) = \hat{l} = -l,$$

where \hat{l} and l are respectively the *u.s.c.* and the *l.s.c.* hulls of l . Thus

$$w^* \text{-Osc}(l|_K)(\gamma) = 2\|l\|.$$

Pick $\alpha > 0$. For any $n \geq 1$, let

$$F_n = \{\gamma \in K : |z_k(\gamma) - z_j(\gamma)| \leq \varepsilon + \alpha \text{ for all } k, j \geq n\}.$$

It follows from the assumption, that

$$K = \bigcup_{n \geq 1} F_n.$$

Hence, by Baire's lemma, there is a non-void w^* -open set $U \subseteq K$ and an $N \geq 1$ such that $U \subseteq F_N$. It follows that, for all $y \in U$, we have

$$|l(y) - z_N(y)| \leq \varepsilon + \alpha,$$

and since z_N is w^* -continuous, this implies that for all $y \in U$ we have

$$\text{Osc}(l|_K)(y) \leq 2(\varepsilon + \alpha).$$

Hence $\text{dist}(l, Z) \leq \varepsilon + \alpha$, which concludes the proof, since $\alpha > 0$ was arbitrary. \square

By Lemmas IV.6 and IV.7, we have

$$\|PF - F\| = \text{dist}(F, L^1) \leq 6\delta\psi(\delta)\|f\|_1$$

for every $F \in L_f$. Once we know this, it follows that one has

$$(4) \quad \|(PW_i - W_i)|_{L^1}\| \leq 6\delta\psi(\delta), \quad i = 1, 2.$$

We define now $U, V : L^1 \rightarrow L^1$ by

$$U = (PW_1)|_{L^1}, \quad V = (PW_2)|_{L^1}.$$

It follows from (1), (2), (3), and (4) that

$$\|I - UV\| \leq \frac{22\delta}{1 - 3\delta} \quad \text{and} \quad \|I - VU\| \leq \frac{22\delta}{1 - 3\delta},$$

so UV and VU , and thus U and V as well, are invertible for $\delta < \frac{1}{25}$. Moreover, $\|U^{-1}\|$ and $\|V^{-1}\|$ are less than $1/(1 - 25\delta)$. We now observe that, since

$$SQ_X f = Q_Y^{**} W_1 f \quad \text{for all } f \in L^1,$$

one has $W_1(X) \subseteq Y^{\perp\perp}$. It follows that $\text{dist}(Uf, Y^{\perp\perp}) \leq \|(PW_1 - W_1)|_{L^1}\|$, so

by (4)

$$\text{dist}(Uf, Y) \leq 2 \text{dist}(Uf, Y^{\perp\perp}) < \frac{13\delta}{1-3\delta},$$

for any $f \in B_X$. It follows that

$$U(B_X) \subseteq B_Y + \frac{27\delta}{1-3\delta} B_{L^1}.$$

Since the same inclusion holds when we exchange X and Y and replace U by V , we can find, for every $g \in B_Y$, an $f \in X$ such that $\|Vg - f\|_1 \leq 13/(1-3\delta)$. Letting $h = f/\|f\|_1$, we obtain

$$\|g - Uh\|_1 \leq \|g - V^{-1}f\|_1 + \|V^{-1}\| \|I - UV\| \|f\|_1 + \|U\|(\|f\|_1 - 1),$$

so

$$B_Y \subseteq U(B_X) + \frac{71\delta}{1-25\delta} B_{L^1},$$

and this finishes the proof of Theorem IV.3. □

Let us say that $Y \subseteq L^1$ has the λ -*lifting property* if for any $T : L^1 \rightarrow L^1/Y$, there exists $\tilde{T} : L^1 \rightarrow L^1$ such that $T = Q_Y \tilde{T}$ and $\|\tilde{T}\| \leq \lambda \|T\|$. If Y is complemented in Y^{**} , it has the lifting property (see [19], Lemma 2.1).

Corollary IV.8. *Let X and Y be small subspaces of L^1 such that there is an isomorphism $S : L^1/X \rightarrow L^1/Y$, such that $\|S\| \cdot \|S^{-1}\| \leq 1 + \delta$, and assume that Y has the λ -lifting property. Then for $\delta \leq \delta(\lambda)$, there is $T : L^1 \rightarrow L^1$, invertible, such that $T(X) = Y$, and*

$$\|T\| \cdot \|T^{-1}\| \leq 1 + K(\lambda)\delta,$$

where $K(\lambda)$ only depends on λ .

Proof. Since $SQ_X = Q_Y^{**}W_{1|L^1}$, by (4) one has, for $\delta \leq \delta_0$,

$$\|SQ_X - Q_Y U\|_{L^1 \rightarrow L^1/Y} \leq C\delta.$$

By assumption, there is $E : L^1 \rightarrow L^1$ with

$$\begin{cases} \|E\| \leq C\lambda\delta, \\ Q_Y E = SQ_X - Q_Y U. \end{cases}$$

Now, if $T = U + E$, one has $SQ_X = Q_Y T$, and the conclusion follows, since this implies, when T is invertible, that $T(X) = Y$. \square

Remark IV.9. There exist subspaces X and Y of infinite codimension in c_0 whose duals are isometric, but for which there is no invertible isometry of ℓ_1 which maps X^\perp onto Y^\perp . Indeed, every isometry from ℓ_1 onto itself is w^* -continuous; hence there should exist an isometry of c_0 mapping X onto Y . Therefore, it suffices to construct infinite codimensional isometric subspaces X and Y of c_0 , in such a way that X contains an element supported by a single integer, and that the elements of Y are all supported by an even number of integers. Since invertible isometries of c_0 leave the size of the supports invariant, such X and Y cannot be mapped into each other by an isometry of all c_0 . So Proposition IV.2 with ℓ_1 instead of L^1 is false in full generality. That shows that it is essential that the measure space be purely non-atomic, at least for defining a proper notion of “small” subspace.

Remark IV.10. It is asked in [19] whether, for two infinite Sidon sets S_1 and S_2 , the spaces $L^1_{S_1}$ and $L^1_{S_2}$ are isomorphic. It would be true if one could lift isomorphisms from $L^1/L^1_{S_1}$ onto $L^1/L^1_{S_2}$. But $L^1_{S_1}$ and $L^1_{S_2}$ are “large” subspaces of L^1 , and the present techniques are apparently not applicable to this question.

V. Construction of a peculiar subspace of L^1 . In Part II, we have seen a unified approach for Alspach’s and Dor’s theorems. In this section, we will see that the results of Part II have limitations. More precisely, we will construct a subspace X of L^1 which will show that the nearness to isometries of the small into-isomorphisms from L^1 into L^1 (and of the small into-isomorphisms from finite dimensional subspaces of L^1 to L^1 ([21])) does not keep on to take place for general subspaces of L^1 . Moreover, the same space will show that extension of isometries from a subspace of L^1 to the whole space L^1 (see [15], [26], [31]) is no longer true if we replace isometries by almost isometries. It shows also that Dor’s Theorem B ([7]), which says that sequences of functions f_n in L^1 which are equivalent to the canonical basis of ℓ_1 are essentially supported by disjoint measurable sets, cannot be improved by replacing the f_n ’s by finite dimensional subspaces: our space X is, for every $\varepsilon > 0$, $(1 + \varepsilon)$ -isomorphic to a ℓ_1 -sum of finite dimensional spaces, but $\|E^S - Id\|_{L(X, L^1)} \geq 1$ for every σ -algebra S generated by a measurable partition. Note also that in our terminology, the space X we construct is small and nicely placed. We work in this section with real Banach spaces.

Theorem V.1. *There exists a subspace X of $L^1 = L^1(\Omega, \Sigma, \mathbb{P})$ such that*

- (a) For every isometry $T : X \rightarrow L^1$ there is a unique isometry $\tilde{T} : L^1 \rightarrow L^1$ with $\tilde{T}|_X = T$;
- (b) For every $\varepsilon > 0$, there exists an isomorphism $T_\varepsilon : X \rightarrow L^1$, with $\|T_\varepsilon\| \|T_\varepsilon^{-1}\| \leq 1 + \varepsilon$, but for which $\|T_\varepsilon - J\| \geq \frac{1}{2}$ for every isometry $J : X \rightarrow L^1$.

Before producing the proof, we note that (b) and Alspach’s theorem (Theorem II.7) show that the operators T_ε do not extend to L^1 with a proper control on the norm of the extension.

Proof. The subspace X is constructed in a very similar way as the one given in [12], Example 4.1.1, but we work with the stable variables themselves instead of working with their absolute value.

The following lemma will be used.

Lemma V.2. Let $p \geq 1$, $\varepsilon > 0$, and $k, l \geq 1$. There exists an $N \geq 1$ such that, for every subspace F of $X = \ell_p^N$ with $\dim F \geq N - k$, there exist l norm one vectors y_1, \dots, y_l of ℓ_p^N , with disjoint support, such that $\text{dist}(y_j, F) \leq \varepsilon$ ($1 \leq j \leq l$).

Proof. Take H bigger than $(1 + 2/\varepsilon)^k$ and $N = Hl$. Since $\dim(X/Z) = k$, indices i_j ($Hh < i_{2h+1} < i_{2h+2} \leq H(h + 1)$, $0 \leq h \leq l$) can be found such that $\text{dist}(e_{2h+2} - e_{2h+1}, F) \leq \varepsilon$, where $(e_n)_{1 \leq n \leq N}$ is the canonical basis of ℓ_p^N . □

Recall first some probabilistic facts (see [3], pp. 60–61). A random variable Z is called p -stable if its characteristic function is $\hat{Z}(t) \stackrel{\text{def}}{=} \mathbb{E}(e^{itZ}) = e^{-c_p|t|^p}$, where c_p is the normalization constant ($\|Z\|_1 = 1$). Now, if $(Z_n)_{n \geq 1}$ is a sequence of independent p -stable variables, its linear span X_p is isometric to ℓ_p , and Z_1, \dots, Z_n, \dots correspond to the usual vector basis of ℓ_p . Since $c_p \xrightarrow{p \rightarrow 1} 0$, for every $\varepsilon > 0$, there is an $\alpha > 0$ such that, if $1 < p \leq \alpha$:

$$d_{\mathbb{P}}(f, 0) \leq \varepsilon \quad \text{for all } f \in B_{X_p},$$

where $d_{\mathbb{P}}$ defines the convergence in probability. By the strong law of large numbers, we also have $\lim_{l \rightarrow +\infty} (1/l) \sum_{j=1}^l |Z_j| = 1$ almost surely.

Now the space X will be the closed linear span of the constant function $\mathbf{1}$ and a sequence of independent p_k -stable variables with $p_k \xrightarrow{k \rightarrow +\infty} 1$. More precisely, we choose inductively:

- (1) a sequence $(p_k)_{k \geq 1}$ converging to 1 and a sequence $(\varepsilon_k)_{k \geq 1}$ converging to 0 such that

$$d_{\mathbb{P}}(f, 0) \leq \varepsilon_k \quad \text{for all } f \in B_{X_{p_k}},$$

where X_{p_k} is the linear span of independent p_k -stable variables; these sequences are chosen so that the subspaces X_1, X_2, \dots constructed in the third step nearly create a ℓ_1 sum. More explicitly, ε_{k+1} is chosen so that the condition

$$d_{\mathbb{P}}(f, 0) \leq \varepsilon_{k+1} \quad \text{for all } f \in B_{X_{p_{k+1}}}$$

implies that, for $g \in X_{k+1}$ and $f \in X_1 + \dots + X_k$, we have

$$\|f + g\|_1 \geq \left(1 - \frac{1}{2^{k+2}}\right) (\|f\|_1 + \|g\|_1);$$

- (2) a sequence of integers $(l_k)_{k \geq 1}$ such that $|(1/l_k) \sum_{j=1}^{l_k} |Y_j| - \mathbf{1}| \leq 1/k$ almost surely for every independent p_k -stable variables Y_1, \dots, Y_{l_k} ;
- (3) a sequence of integers $(N_k)_{k \geq 1}$ such that, if X_k is the linear span of the independent p_k -stable variables $Z_j, j \in I_k$ ($I_k \subseteq \mathbb{N}, |I_k| = N_k$), we can find, by Lemma V.2, for every finite dimensional subspace F of X_k, l_k independent p_k -stable variables Y_1, \dots, Y_{l_k} in X_k such that $\| \mathbf{1} - \text{dist}(Y_j, F) \| \leq 1/2^k$ for $1 \leq j \leq l_k$.

The sets I_1, I_2, \dots are chosen as successive intervals of \mathbb{N} , so we have $\mathbb{N} = \bigcup_{k \geq 1} I_k$. All the p_k -stable variables for the different values of k are independently chosen. Finally, we may and do suppose that the σ -algebra $\sigma(X)$ generated by X is all Σ , since if not we replace $L^1(\Omega, \Sigma, \mathbb{P})$ by $L^1(\Omega, \sigma(X), \mathbb{P}_{\sigma(X)})$. Since $\mathbf{1} \in X$, Hardin's theorem ([15], Corollary 4.3) now gives part (a) of the Theorem.

For the part (b), we state two facts. The first one follows directly from the construction and the almost isometric embedding of finite dimensional subspaces of L^1 into ℓ_1 .

Fact 1. The Banach-Mazur distance from X to the set of subspaces of ℓ_1 is 1.

It follows that for every $\varepsilon > 0$, there exist an operator $T_\varepsilon : X \rightarrow L^1$ such that $\|T_\varepsilon\| \cdot \|T_\varepsilon^{-1}\| \leq 1 + \varepsilon$ and a σ -algebra S_ε generated by disjoint measurable parts of Ω such that $\mathbb{E}^{S_\varepsilon} T_\varepsilon = T_\varepsilon$.

Fact 2. For every sub- σ -algebra S of Σ generated by disjoint measurable sets of Ω , one has

$$\sup_{f \in B_X} \|\mathbb{E}^S f - f\|_1 \geq 1.$$

Proof. Denote by S_1, S_2, \dots a measurable partition of Ω generating S , with $\mathbb{P}(S_i) > 0$ for all i . For every $k \geq 1$, let

$$F_k = \left\{ f \in X_k \mid \int_{S_i} f d\mathbb{P} = 0 \text{ for all } i \leq k \right\}.$$

Let $\varepsilon > 0$. By the third step of the construction, since $\dim(X_k/F_k) \leq k$, and since $\inf_{i \leq k} \mathbb{P}(S_i) > 0$, there is an integer k_0 such that we can find, for every $k \geq k_0$, l_k independent p_k -stable variables Y_1, \dots, Y_{l_k} in X_k so that

$$(*) \quad \left| \frac{1}{\mathbb{P}(S_i)} \int_{S_i} Y_j d\mathbb{P} \right| \leq \varepsilon \text{ for all } j \leq l_k, \text{ all } i \leq k.$$

Let $U_k = S_1 \cup \dots \cup S_k$. By (*), one has

$$\| \mathbb{E}^S(Y_j) \mathbf{1}_{U_k} \|_1 \leq \varepsilon \text{ for all } j \leq l_k.$$

But, by the step 2) of the construction, there is an integer k_1 such that

$$\left\| \frac{1}{l_k} \sum_{j=1}^{l_k} |Y_j| \mathbf{1}_{U_k} \right\|_1 \geq 1 - \varepsilon$$

for $k \geq k_1$, and then there is a $j_k \leq l_k$ such that

$$\| |Y_{j_k}| \mathbf{1}_{U_k} \|_1 \geq 1 - \varepsilon.$$

Hence

$$\| Y_{j_k} - \mathbb{E}^S Y_{j_k} \|_1 \geq \| Y_{j_k} \mathbf{1}_{U_k} - (\mathbb{E}^S Y_{j_k}) \mathbf{1}_{U_k} \|_1 \geq 1 - 2\varepsilon,$$

and this proves the Fact 2. □

More generally, one has the following result:

Lemma V.3. For every isometry $J : X \rightarrow L^1$ and every σ -algebra generated by a measurable partition, one has

$$\sup_{f \in B_X} \| \mathbb{E}^S Jf - Jf \|_1 \geq 1.$$

Proof. From (a), one has an isometry $\tilde{J} : L^1 \rightarrow L^1$ which extends J .

There are measurable functions a, σ on Ω (see Section II) such that

$$\tilde{J}f(s) = a(s)f(\sigma(s)).$$

Since $\|a\|_1 = \|\tilde{J}\mathbf{1}\|_1 = 1$, the measure

$$dQ = |a|d\mathbb{P}$$

is a probability measure. Moreover,

$$\int_{\Omega} |1 + f(t)|d\mathbb{P}(t) = \int_{\Omega} |1 + f(\sigma(s))|dQ(s);$$

hence f and $f(\sigma)$ have the same distribution ([15], Theorem 1.1); so $f(\sigma)$ is p -stable in $L^1(Q)$ whenever f is p -stable in $L^1(\mathbb{P})$. Let S be the σ -algebra generated by a measurable partition (S_i) of Ω , and let S' be a refinement of this partition, so that the sign of a is constant on each S'_i (> 0 , < 0 , or $= 0$). In particular, $\text{supp}(a) = \cup\{S'_i \mid a \neq 0 \text{ on } S'_i\}$.

Let $\varepsilon' > 0$. Reproducing the proof of the above Fact 2 in $L^1(Q)$, we obtain that, if S'_i ($i \leq k$) is disjoint from $\text{supp}(a)$, there exist, for every $k \geq k_0$, $Y_1, \dots, Y_k \in X_k$ such that

$$\left| \frac{1}{Q(S'_i)} \int_{S'_i} Y_j(\sigma) dQ \right| \leq \varepsilon' \quad \text{for all } i, \text{ all } j \leq k,$$

that is

$$\left| \frac{1}{Q(S'_i)} \int_{S'_i} Y_j(\sigma) |a| d\mathbb{P} \right| \leq \varepsilon' \quad \text{for all } i, \text{ all } j \leq k.$$

But a has a constant sign on S'_i , so this writes

$$\left| \frac{1}{Q(S'_i)} \int_{S'_i} Y_j(\sigma) a d\mathbb{P} \right| \leq \varepsilon' \quad \text{for all } i, \text{ all } j \leq k.$$

Therefore

$$\left| \frac{1}{\mathbb{P}(S'_i)} \int_{S'_i} Y_j(\sigma) a d\mathbb{P} \right| \leq \frac{Q(S'_i)}{\mathbb{P}(S'_i)} \varepsilon' \quad \text{for all } i, \text{ all } j \leq k.$$

Consequently, for every $\varepsilon > 0$, there is k'_0 such that, for $k \geq k'_0$, there exist

$Y_1, \dots, Y_{l_k} \in X_k$ such that

$$\left| \frac{1}{\mathbb{P}(S'_i)} \int_{S'_i} Y_j(\sigma) a d\mathbb{P} \right| \leq \varepsilon \quad \text{for all } i, \text{ all } j \leq l_k.$$

Now, since indices $i \leq k$ for which S'_i is disjoint from $\text{supp}(a)$ do not matter, one has, as in the proof of Fact 2, for k with $U_k = S'_1 \cup \dots \cup S'_{l_k} \in S$:

$$\|\mathbb{E}^S J(Y_j) \mathbf{1}_{U_k}\|_1 \leq \|\mathbb{E}^{S'} J(Y_j) \mathbf{1}_{U_k}\|_1 \leq \varepsilon \quad \text{for all } j \leq l_k.$$

On the other hand,

$$\frac{1}{l_k} \sum_{j=1}^{l_k} |Y_j(\sigma)| \xrightarrow[k \rightarrow +\infty]{} \mathbf{1} \quad \mathbb{P}\text{-a.s.},$$

$$\text{hence} \quad \frac{1}{l_k} \sum_{j=1}^{l_k} |J(Y_j)| \xrightarrow[k \rightarrow +\infty]{} |a| \quad Q\text{-a.s.}$$

Since $\|a\|_{L^1(\mathbb{P})} = 1$, it follows that

$$\varliminf_{k \rightarrow +\infty} \left\| \frac{1}{l_k} \sum_{j=1}^{l_k} |J(Y_j)| \mathbf{1}_{U_k} \right\|_1 \geq 1.$$

There is therefore a k_1 such that, for every $k \geq k_1$, there is a $j_k \leq l_k$ such that

$$\|J(Y_{j_k}) \mathbf{1}_{U_k}\|_1 \geq 1 - \varepsilon,$$

and so

$$\|J(Y_{j_k}) - \mathbb{E}^S J(Y_{j_k})\|_1 \geq \|J(Y_{j_k}) \mathbf{1}_{U_k} - \mathbb{E}^S J(Y_{j_k}) \mathbf{1}_{U_k}\|_1 \geq 1 - 2\varepsilon,$$

and this ends the proof of the lemma. □

We can now finish the proof of Theorem V.1.

Let $\varepsilon > 0$, $T_\varepsilon : X \rightarrow L^1$ with $\|T_\varepsilon\| \cdot \|T_\varepsilon^{-1}\| \leq 1 + \varepsilon$, and S_ε such that $\mathbb{E}^{S_\varepsilon} T_\varepsilon = T_\varepsilon$, as previously defined. For every isometry $J : X \rightarrow L^1$, we have

$$\|\mathbb{E}^{S_\varepsilon} J - T_\varepsilon\| = \|\mathbb{E}^{S_\varepsilon} J - \mathbb{E}^{S_\varepsilon} T_\varepsilon\| \leq \|J - T_\varepsilon\|,$$

so

$$\|\mathbb{E}^{S_\varepsilon} J - J\| \leq \|\mathbb{E}^{S_\varepsilon} J - T_\varepsilon\| + \|T_\varepsilon - J\| \leq 2\|T_\varepsilon - J\|.$$

But from Lemma V.3, $\|\mathbb{E}^{S_\varepsilon} J - J\| \geq 1$, and so $\|T_\varepsilon - J\| \geq 1/2$. □

Remarks.

- (1) There is a gap in the proof of the implication (iii)⇒(iv) of Corollary 3.5 in [12], and we do not know whether one has (i)⇒(iv) in that Corollary 3.5. Anyway, Theorem V.1 shows that one cannot replace the distance d_m in [12], Corollary 3.5 (iv), by the $\|\cdot\|_1$ distance.
- (2) We mention here that it follows from L. Schwartz’s thesis ([34]; see [2], Theorem 4.2.5), that Müntz spaces $M_1(\{t^{n_k}\})$ (with $\sum n_k^{-1} < +\infty$) are examples of subspaces of L^1 almost isometric to subspaces of ℓ_1 . The case $p > 1$ has been noticed in [5] (see comments after Corollary 1.8).

The derivation of Lemma V.3 from Lemma V.2 can actually be put in a more general frame. This is the content of the next proposition.

Proposition V.4. *Let X_1 and X_2 be isometric subspaces of $L^1(\Omega_1, \Sigma_1, \mu_1)$ and $L^1(\Omega_2, \Sigma_2, \mu_2)$ respectively. Suppose that X_1 contains the constant functions, and that, for some $\varepsilon > 0$, there is a σ -algebra \mathcal{A} , generated by a measurable partition of Ω_1 , for which*

$$\|f - \mathbb{E}^{\mathcal{A}} f\|_1 \leq \varepsilon \|f\|_1 \quad \text{for all } f \in X_1.$$

Then, for $\varepsilon' > \varepsilon$, there is a $(\Sigma_2 \otimes \text{Bor})$ -measurable partition of $\Omega_2 \times [0, 1]$, generating a σ -algebra \mathcal{B} , for which one has

$$\|g - \mathbb{E}^{\mathcal{B}} g\|_1 \leq \varepsilon' \|g\|_1 \quad \text{for all } g \in X_2.$$

In this statement, we identify $L^1(\Omega_2, \Sigma_2, \mu_2)$ to a subspace of

$$L^1(\Omega_2 \times [0, 1], \Sigma_2 \otimes \text{Bor}, \mu_2 \otimes dt) = L^1(\tilde{\Omega}_2, \tilde{\Sigma}_2, \tilde{\mu}_2),$$

by letting $\tilde{g}(\omega, t) = g(\omega)$ for $\omega \in \Omega_2$, $t \in [0, 1]$, and $g \in L^1(\Omega_2, \Sigma_2, \mu_2)$.

Proof. Denote by A_1, A_2, \dots the partition generating \mathcal{A} , and by σ_1 and σ_2 the sub- σ -algebras of Σ_1 and Σ_2 generated by X_1 and X_2 respectively. Let U be an isometry from X_1 onto X_2 . Since $\mathbf{1} \in X_1$, Hardin’s theorem ([15], Corollary 4.3) ensures the existence of an into-isometry

$$T = \tilde{U} : L^1(\sigma_1) \rightarrow L^1(\Sigma_2)$$

which extends U , and whose range is $L^1(\sigma_2)$. T can be written as $Tf = a(f \circ \tau)$, with τ an isomorphism between (Ω_2, σ_2) and (Ω_1, σ_1) . Write $Jf = f \circ \tau$. J is positive, and so, since $\varphi_n = \mathbb{E}^{\sigma_1}(\mathbf{1}_{A_n})$ verifies

$$\varphi_n \geq 0 \quad \text{and} \quad \sum_{n \geq 1} \varphi_n = \mathbf{1},$$

we have $\psi_n = J\varphi_n \geq 0$ and $\sum_{n \geq 1} \psi_n = \mathbf{1}$. Note that $\nu = |a|\mu_2$ is a probability measure and J is an isometry from $L^1(\Omega_1, \sigma_1, \mu_1)$ to $L^1(\Omega_2, \sigma_2, \nu)$. We have:

Lemma V.5. *If ψ_n are positive σ_2 -measurable functions with $\sum_{n \geq 1} \psi_n = \mathbf{1}$, there is a measurable partition in sets $B_1, B_2, \dots \in \sigma_2 \otimes \text{Bor}$, such that $\tilde{\psi}_n = \mathbb{E}^{\sigma_2}(\mathbf{1}_{B_n})$, for $n \geq 1$.*

Proof. Set

$$B_n = \left\{ (\omega, t) \mid \sum_{k=1}^{n-1} \psi_k(\omega) \leq t < \sum_{k=1}^n \psi_k(\omega) \right\}.$$

We have for every $B \in \sigma_2$:

$$\begin{aligned} \int_{\tilde{B}} \mathbb{E}^{\sigma_2} \mathbf{1}_{B_n} &= \int_{\tilde{B}} \mathbf{1}_{B_n} = (\nu \otimes dt)(B_n \cap \tilde{B}) \\ &= \int_B m[B_n(\omega)] d\nu(\omega) = \int_B \psi_n(\omega) d\nu(\omega), \end{aligned}$$

so $\tilde{\psi}_n = \mathbb{E}^{\sigma_2} \mathbf{1}_{B_n}$. □

We may suppose now that \tilde{a} has a constant sign on each B_n . Indeed, since \tilde{a} is σ_2 -measurable and τ is an isomorphism from (Ω_2, σ_2) onto (Ω_1, σ_1) , we have, if $B_n^+ = B_n \cap \{\tilde{a} \geq 0\}$ and $B_n^- = B_n \cap \{\tilde{a} < 0\}$,

$$\begin{aligned} \mathbb{E}^{\sigma_2}(\mathbf{1}_{B_n^+}) &= \mathbf{1}_{\{\tilde{a} \geq 0\}} \mathbb{E}^{\sigma_2}(\mathbf{1}_{B_n}) = \mathbf{1}_{\{\tilde{a} \geq 0\}} \tilde{\psi}_n \\ &= \mathbf{1}_{\{\tilde{a} \geq 0\}} (\tilde{\varphi}_n \circ \tau) = \tilde{\mathbb{E}}^{\sigma_1}(\mathbf{1}_{\{a \circ \tau^{-1} \geq 0\}} \mathbf{1}_{A_n}) \circ \tau, \end{aligned}$$

so we may cut each A_n with the sets $\{a \circ \tau^{-1} \geq 0\}$ and $\{a \circ \tau^{-1} < 0\}$, and the B_n^+ 's and the B_n^- 's will correspond to the A_n^+ 's and the A_n^- 's. Now, letting

$$a_n(f) = \frac{1}{\mu_1(A_n)} \int_{A_n} f d\mu_1,$$

one has

$$\begin{aligned}
 \text{(i)} \quad \mu_1(A_n) &= \int_{\Omega_1} \mathbf{1}_{A_n} d\mu_1 = \int_{\Omega_1} \varphi_n d\mu_1 = \|\varphi_n\|_1 = \|\psi_n\|_{L^1(\nu)} \\
 &= \int_{\Omega_2} \psi_n d\nu = \int_{\tilde{\Omega}_2} \mathbf{1}_{B_n} d\tilde{\nu} = \tilde{\nu}(B_n),
 \end{aligned}$$

and

$$\begin{aligned}
 \text{(ii)} \quad \int_{A_n} f d\mu_1 &= \int_{\Omega_1} f \mathbf{1}_{A_n} d\mu_1 = \int_{\Omega_1} f \varphi_n d\mu_1 = \int_{\Omega_2} J(f \varphi_n) d\nu \\
 &\qquad \text{since } Jf = f \circ \tau, \text{ and } \mu_1 = \tau^*(\nu), \\
 &= \int_{\Omega_2} (Jf) \psi_n d\nu = \int_{B_n} \tilde{J}f d\tilde{\nu}.
 \end{aligned}$$

Hence

$$a_n(f) = \frac{1}{\tilde{\nu}(B_n)} \int_{B_n} \tilde{J}f d\tilde{\nu},$$

and then, denoting by \mathcal{B} the σ -algebra generated by B_1, B_2, \dots ,

$$\begin{aligned}
 \|\tilde{J}f - \mathbb{E}^{\mathcal{B}}(\tilde{J}f)\|_{L^1(\tilde{\nu})} &= \sum_{n \geq 1} \int_{B_n} |\tilde{J}f - a_n(f) \mathbf{1}| d\tilde{\nu} = \sum_{n \geq 1} \int_{B_n} |f \circ \tau - a_n(f) \mathbf{1}| d\tilde{\nu} \\
 &= \sum_{n \geq 1} \int_{\Omega_2} |f \circ \tau - a_n(f) \mathbf{1}| \psi_n d\nu \\
 &= \sum_{n \geq 1} \int_{\Omega_2} |f \circ \tau - a_n(f) \mathbf{1}| (\varphi_n \circ \tau) d\nu \\
 &= \sum_{n \geq 1} \int_{\Omega_1} |f - a_n(f) \mathbf{1}| \varphi_n d\mu_1 = \sum_{n \geq 1} \int_{A_n} |f - a_n(f) \mathbf{1}| d\mu_1 \\
 &= \|f - \mathbb{E}^{\mathcal{A}} f\|_1 \leq \varepsilon \|f\|_1.
 \end{aligned}$$

Now, going back to T , we refine \mathcal{B} in such a way that $\|\tilde{a} - \mathbb{E}^{\mathcal{B}} \tilde{a}\|_{L^1(\tilde{\mu}_2)} \leq \alpha = \varepsilon' - \varepsilon$. We keep the inequality $\|\tilde{J}f - \mathbb{E}^{\mathcal{B}}(\tilde{J}f)\|_{L^1(\tilde{\nu})} \leq \varepsilon \|f\|_1$. Moreover, this refinement can be made by taking the intersection with σ_2 -measurable sets, so, as above, the new B_n 's will correspond to an appropriate refinement of the A_n 's. Set $\varepsilon_n = 1$ if $a \geq 0$ on B_n , and $\varepsilon_n = -1$ if $a < 0$ on B_n . One has

$$\int_{B_n} \tilde{a} d\tilde{\mu}_2 = \varepsilon_n \int_{\Omega_2} |a| \psi_n d\mu_2 = \varepsilon_n \int |a| (\varphi_n \circ \tau) d\mu_2 = \varepsilon_n \int_{\Omega_2} |T\varphi_n| d\mu_2$$

$$= \varepsilon_n \|T\varphi_n\|_{L^1(\mu_2)} = \varepsilon_n \|\varphi_n\|_{L^1(\mu_1)} = \varepsilon_n \mu_1(A_n),$$

and so

$$\mathbb{E}^B \tilde{a} = \sum_{n \geq 1} \varepsilon_n \frac{\mu_1(A_n)}{\tilde{\mu}_2(B_n)} \mathbf{1}_{B_n}.$$

Now

$$\begin{aligned} & \|\tilde{T}f - \mathbb{E}^B \tilde{T}f\|_{L^1(\tilde{\mu}_2)} \\ &= \int_{\tilde{\Omega}_2} \left| \tilde{a}(\tilde{f} \circ \tau) - \sum_{n \geq 1} \left(\frac{1}{\tilde{\mu}_2(B_n)} \int_{B_n} \tilde{a}(\tilde{f} \circ \tau) \right) \mathbf{1}_{B_n} \right| d\tilde{\mu}_2 \\ &\leq \int_{\tilde{\Omega}_2} \left| \tilde{a}(\tilde{f} \circ \tau) - \sum_{n \geq 1} \left(\frac{\tilde{a}}{\mu_1(A_n)} \int_{B_n} |\tilde{a}|(\tilde{f} \circ \tau) d\tilde{\mu}_2 \right) \mathbf{1}_{B_n} \right| d\tilde{\mu}_2 \\ &\quad + \int_{\tilde{\Omega}_2} \left| \sum_{n \geq 1} \left(\frac{\tilde{a}}{\mu_1(A_n)} \int_{B_n} |\tilde{a}|(\tilde{f} \circ \tau) d\tilde{\mu}_2 \right) \mathbf{1}_{B_n} \right. \\ &\quad \quad \quad \left. - \sum_{n \geq 1} \left(\frac{1}{\tilde{\mu}_2(B_n)} \int_{B_n} \tilde{a}(\tilde{f} \circ \tau) d\tilde{\mu}_2 \right) \mathbf{1}_{B_n} \right| d\tilde{\mu}_2 \\ &\leq \int_{\tilde{\Omega}_2} \left| \tilde{f} \circ \tau - \sum_{n \geq 1} \left(\frac{1}{\tilde{\nu}(B_n)} \int_{B_n} (\tilde{f} \circ \tau) d\tilde{\nu} \right) \mathbf{1}_{B_n} \right| d\tilde{\nu} \\ &\quad + \int_{\tilde{\Omega}_2} \sum_{n \geq 1} \left(\frac{1}{\mu_1(A_n)} \int_{B_n} |\tilde{a}|(\tilde{f} \circ \tau) d\tilde{\mu}_2 \right) \left(\tilde{a} - \varepsilon_n \frac{\mu_1(A_n)}{\tilde{\mu}_2(B_n)} \right) \mathbf{1}_{B_n} d\tilde{\mu}_2 \\ &\leq \|\tilde{J}f - \mathbb{E}^B \tilde{J}f\|_{L^1(\tilde{\nu})} + \int_{\tilde{\Omega}_2} \sum_{n \geq 1} \|f\|_1 \left| \tilde{a} - \varepsilon_n \frac{\mu_1(A_n)}{\tilde{\mu}_2(B_n)} \right| \mathbf{1}_{B_n} d\tilde{\mu}_2 \\ &\leq \varepsilon \|f\|_1 + \|f\|_1 \sum_{n \geq 1} \int_{B_n} \left| \tilde{a} - \varepsilon_n \frac{\mu_1(A_n)}{\tilde{\mu}_2(B_n)} \right| d\tilde{\mu}_2 \\ &\leq \varepsilon \|f\|_1 + \|f\|_1 \|\tilde{a} - \mathbb{E}^B \tilde{a}\|_{L^1(\tilde{\mu}_2)} \leq (\varepsilon + \alpha) \|f\|_1 = \varepsilon' \|f\|_1. \quad \square \end{aligned}$$

Remarks.

- (1) A slight variation in Proposition V.4 is: set $C_n = \tau^{-1}(A_n)$. The C_n 's are in general not in Σ_2 , but if C is the σ -algebra generated by them, the map $\tau : \Omega_2 \rightarrow \Omega_1$ is $(\sigma_2 \vee C - \sigma_1 \vee \mathcal{A})$ -bi-measurable (see in [14] a description of these σ -algebras), and $\tilde{\mu}_2(B) = \mu_1[\tau(B)]$ defines a measure on $(\Omega_2, \sigma_2 \vee C)$, whose restriction to σ_2 is equal to the restriction of μ_2 to σ_2 . One has then an isometry $\tilde{J} : L^1(\Omega_1, \sigma_1 \vee \mathcal{A}, \mu_1) \rightarrow L^1(\Omega_2, \sigma_2 \vee C, \tilde{\mu}_2)$, defined by $\tilde{J}(f) = f \circ \tau$, such that $\tilde{J}X_1 = JX_1$. It follows that $\|Jf - \mathbb{E}^C Jf\|_1 \leq \varepsilon \|f\|_1$. It would be interesting to find an

intrinsic characterization of these spaces. This problem may be connected to the fact that, though there are sufficient conditions to have uniform convergence of martingales (see [29]), it seems there is at the present time no available necessary condition.

- (2) Since the space X and $T_\varepsilon(X)$ are small nicely placed subspaces (by Proposition III.7), it follows from Theorem V.1 and Proposition IV.2 that there is $\alpha > 0$ such that the quotient spaces L^1/X and $L^1/T_\varepsilon(X)$ have Banach-Mazur distance greater than $(1 + \alpha)$ for all $\varepsilon > 0$. This implies of course that the operators T_ε do not extend to isomorphisms U of L^1 such that the norms of U and its inverse are small.

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G. GODEFROY

Université Pierre et Marie Curie

Equipe d'Analyse, Tour 46-0, Case 186, 4 place Jussieu

F-75252 Paris Cedex 05, FRANCE

EMAIL: gig@ccr.jussieu.fr

N. J. KALTON

Department of Mathematics
University of Missouri-Columbia
Columbia, MO 65211 U. S. A.
EMAIL: nigel@math.missouri.edu

D. LI

Faculté des Sciences Jean Perrin
Université d'Artois, Pôle de Lens
Rue Jean Souvraz, SP 18
F-62307 Lens Cedex, FRANCE
EMAIL: daniel.li@euler.univ-artois.fr

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