# The One-Third-Trick and Shift Operators by <br> Richard LECHNER <br> Presented by Gilles PISIER 

Summary. We obtain a representation as martingale transform operators for the rearrangement and shift operators introduced by T. Figiel. The martingale transforms and the underlying sigma algebras are obtained explicitly by combinatorial means. The known norm estimates for those operators are a direct consequence of our representation.

1. Introduction. The proof of the $T(1)$ theorem by T. Figiel proceeds by expanding the integral operator into an absolutely converging series of basic building blocks $T_{m}$ and $U_{m}$, rearranging and shifting the Haar system. This involves the following norm estimates for those building blocks, which T. Figiel obtained by combinatorial means:

$$
\begin{align*}
& \left\|T_{m}: L_{X}^{p} \rightarrow L_{X}^{p}\right\| \leq C\left(\log _{2}(2+|m|)\right)^{\alpha},  \tag{1.1}\\
& \left\|U_{m}: L_{X}^{p} \rightarrow L_{X}^{p}\right\| \leq C\left(\log _{2}(2+|m|)\right)^{\beta}, \tag{1.2}
\end{align*}
$$

where the constant $C>0$ depends only on $p$, the UMD-constant of $X$ and $0<\alpha, \beta<1$. For the original proof see [Fig88] and [FW01]. See also [NS97] and Mül05. For extensions to spaces of homogeneous type see MP12].

The purpose of the present paper is to obtain a representation of $T_{m}$ and $U_{m}$ as the sum of roughly $\log _{2}(2+|m|)$ martingale transform operators. This is done by combinatorial analysis of the equations defining $T_{m}$ and $U_{m}$.

Our combinatorial analysis exhibits the link of T. Figiel's rearrangement and shift operators to the so called one-third-trick originating in the work of Wol82, [GJ82], Dav80] and CWW85.

[^0]Related recent developments. Recently T. P. Hytönen Hyt12 presented his own proof of T. Figiel's vector-valued $T(1)$ theorem (see [Fig90). The basic aim of T. P. Hytönen in Hyt12 is the same as ours, to find reductions of the general case to certain preferable situations where the so called Figiel compatibility condition is satisfied. To this end the problem in Hyt12 is randomized and the properties of the so called random dyadic partitions of Nazarov, Treil and Volberg [NTV97, NTV03] are exploited.

By contrast, the proof in the present paper proceeds by finding explicitly those filtrations that turn a given bad combinatorial situation into a good one, such that Figiel's compatibility condition is satisfied. Our reduction is self-contained and develops specific combinatorics of colored dyadic intervals.

## 2. Preliminaries

The Haar system in $\mathbb{R}$. For any given finite interval $I \subset \mathbb{R}$ we define

$$
h_{I}=1_{I_{0}}-1_{I_{1}},
$$

where $1_{A}$ denotes the characteristic function of a set $A, I_{0}=[\inf I,(\inf I+$ $\sup I) / 2\left[\right.$ and $I_{1}=\left[(\inf I+\sup I) / 2, \sup I\left[\right.\right.$ Let $\mathscr{D}_{j}=\left\{\left[2^{-j} k, 2^{-j}(k+1)[:\right.\right.$ $k \in \mathbb{Z}\}, j \in \mathbb{Z}$, and define the collection of dyadic intervals by $\mathscr{D}=\bigcup_{j \in \mathbb{Z}} \mathscr{D}_{j}$. The $L^{\infty}$-normalized Haar system is given by $\left\{h_{I}: I \in \mathscr{D}\right\}$.

Banach spaces with the UMD-property. By $L^{p}(\Omega, \mu ; X)$ we denote the space of functions with values in $X$, Bochner-integrable with respect to $\mu$. If $\Omega=\mathbb{R}$ and $\mu$ is the Lebesgue measure $|\cdot|$ on $\mathbb{R}$, then set $L_{X}^{p}(\mathbb{R})=$ $L^{p}(\mathbb{R},|\cdot| ; X)$, if unambiguous further abbreviated as $L_{X}^{p}$.

We say $X$ is a $U M D$ space if for every $X$-valued martingale difference sequence $\left\{d_{j}\right\}_{j} \subset L^{p}(\Omega, \mu ; X), 1<p<\infty$, and choice of signs $\varepsilon_{j} \in\{-1,1\}$ one has

$$
\begin{equation*}
\left\|\sum_{j} \varepsilon_{j} d_{j}\right\|_{L^{p}(\Omega, \mu ; X)} \leq \mathscr{U}_{p}(X)\left\|\sum_{j} d_{j}\right\|_{L^{p}(\Omega, \mu ; X)}, \tag{2.1}
\end{equation*}
$$

where $\mathscr{U}_{p}(X)$ does not depend on $\varepsilon_{j}$ or $d_{j}$. The constant $\mathscr{U}_{p}(X)$ is called the UMD-constant. We refer the reader to [Bur81].

Kahane's contraction principle. For every Banach space $X, 1 \leq p$ $<\infty$, finite set $\left\{x_{j}\right\}_{j} \subset X$ and bounded sequence $\left\{c_{j}\right\}_{j}$ of scalars we have

$$
\begin{equation*}
\left(\int_{0}^{1}\left\|\sum_{j} r_{j}(t) c_{j} x_{j}\right\|_{X}^{p} d t\right)^{1 / p} \leq \sup _{j}\left|c_{j}\right|\left(\int_{0}^{1}\left\|\sum_{j} r_{j}(t) x_{j}\right\|_{X}^{p} d t\right)^{1 / p}, \tag{2.2}
\end{equation*}
$$

where $\left\{r_{j}\right\}_{j}$ denotes an independent sequence of Rademacher functions. For details see Kah85.

Kahane's inequality. Given $1 \leq p<\infty$, there exists a constant $K_{p}$ such that for any Banach space $X$ and any finite sequence $\left\{x_{j}\right\} \subset X$ we have

$$
\begin{equation*}
\left(\int_{0}^{1}\left\|\sum_{j} r_{j}(t) x_{j}\right\|_{X}^{p} d t\right)^{1 / p} \leq K_{p} \int_{0}^{1}\left\|\sum_{j} r_{j}(t) x_{j}\right\|_{X} d t \tag{2.3}
\end{equation*}
$$

where $\left\{r_{j}\right\}_{j}$ denotes an independent sequence of Rademacher functions. For details see Kah85.

The martingale inequality of Stein - Bourgain's version. Let $(\Omega, \mathcal{F}, \mu)$ be a probability space, and let $\mathcal{F}_{1} \subset \cdots \subset \mathcal{F}_{m} \subset \mathcal{F}$ denote an increasing sequence of $\sigma$-algebras. For every choice of $f_{1}, \ldots, f_{m} \in L^{p}(\Omega, \mu ; X)$, $1<p<\infty$, letting $r_{1}, \ldots, r_{m}$ denote independent Rademacher functions, we have

$$
\begin{equation*}
\int_{0}^{1}\left\|\sum_{i=1}^{m} r_{i}(t) \mathbb{E}\left(f_{i} \mid \mathcal{F}_{i}\right)\right\|_{L^{p}(\Omega, \mu ; X)} d t \leq C \int_{0}^{1}\left\|\sum_{i=1}^{m} r_{i}(t) f_{i}\right\|_{L^{p}(\Omega, \mu ; X)} d t \tag{2.4}
\end{equation*}
$$

where $C$ depends only on $p$ and $\mathscr{U}_{p}(X)$.
The scalar-valued version of (2.4) by E. M. Stein can be found in Ste70]. The vector-valued extension is due to J. Bourgain Bou86]. A proof may be found in FW01.
3. The one-third-trick. The one-third-trick originates in the work of Wol82, GJ82] and CWW85.
3.1. Bilateral alternating one-third-shift. For every $j \in \mathbb{Z}$ let

$$
\begin{equation*}
s_{j}=(-1)^{j} 2^{-j} / 3 \tag{3.1}
\end{equation*}
$$

and define

$$
\begin{equation*}
s(I)=s_{j} \tag{3.2}
\end{equation*}
$$

for all intervals $I$ having measure $|I|=2^{-j}$. Then define the one-third-shift map

$$
\begin{equation*}
\sigma(I)=I+s(I) \tag{3.3}
\end{equation*}
$$

and the one-third-shift operator

$$
\begin{equation*}
S\left(h_{I}\right)=h_{\sigma(I)} \tag{3.4}
\end{equation*}
$$

For every given interval $I \in \mathscr{D}$ there exists a unique one-third-shifted interval $J \in \sigma(\mathscr{D})$ with $|J|=|I| / 2$, contained in $I$. Observe that for every $j \in \mathbb{Z}$ and $I \in \mathscr{D}_{j}$ we have

$$
\begin{array}{r}
\#\left\{J \in \sigma\left(\mathscr{D}_{j+1}\right): J \cap I \neq \emptyset\right\}=3 \\
\#\left\{J \in \sigma\left(\mathscr{D}_{j+1}\right): J \subset I\right\}=1
\end{array}
$$

So we can define $\omega(I)$ by

$$
\begin{equation*}
\omega(I)=J, \quad \text { where } \quad J \in \sigma(\mathscr{D}),|J|=|I| / 2 \text { and } J \subset I . \tag{3.5}
\end{equation*}
$$

Note the basic properties summarized in:
Lemma 3.1. The following statements are true:
(i) $\sigma(\mathscr{D})$ is a nested collection of dyadic intervals, and $\mathscr{D} \cap \sigma(\mathscr{D})=\emptyset$.
(ii) $\omega: \mathscr{D} \rightarrow \sigma(\mathscr{D})$ is well defined and injective.
(iii) If $I \in \mathscr{D}$, then $\omega(I) \subset I$.
(iv) For every $I \in \mathscr{D}$ we have $\operatorname{dist}\left(\omega(I), I^{c}\right)=|I| / 6$.
(v) If $I, J \in \mathscr{D}$ and $|I|=|J|$, then $\operatorname{dist}(\omega(I), \omega(J))<|\omega(I)|$ if and only if $I=J$.
(vi) For all $I \in \mathscr{D}$ we have the identity

$$
\sigma(I)=\omega(I) \cup(\omega(I)+\operatorname{sign}(s(I)) \cdot|\omega(I)|)
$$

Proof. The assertions are easily verified.
We need some more notation. For all $j \in \mathbb{Z}$ and

$$
u=\sum_{I \in \mathscr{D}} u_{I} h_{I}
$$

let $(u)_{j}$ restrict the function $u$ to level $j$, precisely

$$
\begin{equation*}
(u)_{j}=\sum_{I \in \mathscr{D}_{j}} u_{I} h_{I} \tag{3.6}
\end{equation*}
$$

If we define

$$
\begin{equation*}
\mathbb{I}(u)_{j}=\sum_{I \in \mathscr{D}_{j}} u_{I} 1_{I} \tag{3.7}
\end{equation*}
$$

then Kahane's contraction principle (2.2) and Kahane's inequality (2.3) yield

$$
\begin{equation*}
\int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t)(u)_{j}\right\|_{L_{X}^{p}} d t \approx \int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t) \mathbb{I}(u)_{j}\right\|_{L_{X}^{p}} d t . \tag{3.8}
\end{equation*}
$$

The following theorem establishes that the one-third-shift operator $S$ : $L_{X}^{p} \rightarrow L_{X}^{p}$ is an isomorphism.

Theorem 3.2. Let $1<p<\infty$ and $X$ a Banach space with the UMDproperty. Then there exists a constant $C>0$ such that

$$
\frac{1}{C}\|u\|_{L_{X}^{p}} \leq\|S u\|_{L_{X}^{p}} \leq C\|u\|_{L_{X}^{p}}
$$

for all $u \in L_{X}^{p}$. The constant $C$ depends only on $\mathscr{U}_{p}(X)$.

Proof. Let $u=\sum_{I \in \mathscr{D}} u_{I} h_{I}|I|^{-1} \in L_{X}^{p}$ be fixed throughout this proof and set

$$
\begin{aligned}
v & =\sum_{I \in \mathscr{D}} u_{I} h_{\omega(I)}|\omega(I)|^{-1}, \\
(v)_{j} & =\sum_{|\omega(I)|=2^{-j}} u_{I} h_{\omega(I)}|\omega(I)|^{-1} .
\end{aligned}
$$

Note that $\{\omega(I): I \in \mathscr{D}\}$ is nested (see Lemma 3.1(ii) \& (iii)). Observe that

$$
\mathbb{I}(u)_{j}=\mathbb{E}\left(\mathbb{I}(v)_{j} \mid \mathscr{D}_{j}\right)
$$

by Lemma 3.1(iii), so the UMD-property, Kahane's contraction principle (2.2) and Kahane's inequality (2.3) yield

$$
\begin{aligned}
\|u\|_{L_{X}^{p}} & \lesssim \int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t)(u)_{j}\right\|_{L_{X}^{p}} d t \\
& =\int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t) \mathbb{I}(u)_{j}\right\|_{L_{X}^{p}} d t \\
& =\int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t) \mathbb{E}\left(\mathbb{I}(v)_{j} \mid \mathscr{D}_{j}\right)\right\|_{L_{X}^{p}} d t
\end{aligned}
$$

Now we apply Stein's martingale inequality (2.4) followed by (3.8) to pass from $\mathbb{I}(v)_{j}$ to $(v)_{j}$, so

$$
\begin{aligned}
\int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t) \mathbb{E}\left(\mathbb{I}(v)_{j} \mid \mathscr{D}_{j}\right)\right\|_{L_{X}^{p}} d t & \lesssim \int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t) \mathbb{I}(v)_{j}\right\|_{L_{X}^{p}} d t \\
& \approx \int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t)(v)_{j}\right\|_{L_{X}^{p}} d t .
\end{aligned}
$$

Recalling definition (3.4) and applying Kahane's contraction principle, in view of $\omega(I) \subset \sigma(I)$ (see identity Lemma 3.1(vi)), we estimate

$$
\int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t)(v)_{j}\right\|_{L_{X}^{p}} d t \leq 2 \int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t)(S u)_{j}\right\|_{L_{X}^{p}} d t
$$

and the UMD-property implies

$$
\int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t)(v)_{j}\right\|_{L_{X}^{p}} d t \lesssim\|S u\|_{L_{X}^{p}}
$$

Thus, collecting the inequalities yields

$$
\|u\|_{L_{X}^{p}} \lesssim\|S u\|_{L_{X}^{p}}
$$

One can repeat the preceding argument with the roles of $u$ and $S u$ interchanged to obtain the converse inequality

$$
\|S u\|_{L_{X}^{p}} \lesssim\|u\|_{L_{X}^{p}} .
$$

3.2. Unilateral one-third-shift. We now introduce modified versions $\sigma_{0}$ and $\sigma_{1}$ of the one-third-shift map $\sigma$. To this end we define $\sigma_{0}, \sigma_{1}: \mathscr{D} \rightarrow$ $\sigma(\mathscr{D})$ by

$$
\begin{array}{lll}
\sigma_{0}(I)=J, & \text { where } & J \in \sigma(\mathscr{D}),|J|=|I| \text { and } \sup J \in I, \\
\sigma_{1}(I)=J, & \text { where } & J \in \sigma(\mathscr{D}),|J|=|I| \text { and } \inf J \in I . \tag{3.10}
\end{array}
$$

This induces the one-third-shift operators $S_{0}$ and $S_{1}$ given by the linear extension of

$$
\begin{array}{ll}
S_{0}\left(h_{I}\right) h_{\sigma_{0}(I)}, & I \in \mathscr{D}, \\
S_{1}\left(h_{I}\right)=h_{\sigma_{1}(I)}, & I \in \mathscr{D} . \tag{3.12}
\end{array}
$$

Observe that if $I \in \mathscr{D}_{j}$ and $j$ is odd we have

$$
\sigma_{0}(I)=\sigma(I) \quad \text { and } \quad \sigma_{1}(I)=\sigma(I)+|I|,
$$

and if $j$ is even

$$
\sigma_{0}(I)=\sigma(I)-|I| \quad \text { and } \quad \sigma_{1}(I)=\sigma(I) .
$$

Furthermore, since

$$
\sigma_{i}\left(\mathscr{D}_{j}\right)=\sigma\left(\mathscr{D}_{j}\right) \quad \text { for all } i=0,1 \text { and } j \in \mathbb{Z},
$$

we deduce from Lemma 3.1(ii) that $\sigma_{i}(\mathscr{D})$ is nested.
Finally, we can see that

$$
\left|I \cap \sigma_{0}(I)\right| \geq \frac{1}{3}|I|, \quad\left|I \cap \sigma_{1}(I)\right| \geq \frac{1}{3}|I|,
$$

for all $I \in \mathscr{D}$. The proof of Theorem 3.2 with minor modifications yields Theorem 3.3 below.

Theorem 3.3. Let $1<p<\infty$ and $X$ a Banach space with the UMDproperty. Then there exists a constant $C>0$ such that

$$
\begin{aligned}
\frac{1}{C}\|u\|_{L_{X}^{p}} & \leq\left\|S_{0} u\right\|_{L_{X}^{p}} \leq C\|u\|_{L_{X}^{p}} \\
\frac{1}{C}\|u\|_{L_{X}^{p}} & \leq\left\|S_{1} u\right\|_{L_{X}^{p}} \leq C\|u\|_{L_{X}^{p}}
\end{aligned}
$$

for all $u \in L_{X}^{p}$. The constant $C$ depends only on $\mathscr{U}_{p}(X)$.
Proof. Define $\omega_{0}$ and $\omega_{1}$ by
$\omega_{0}(I)=J, \quad$ where $\quad J \in \sigma(\mathscr{D}),|J|=|I| / 4$ and $\sup J=\sup \sigma_{0}(I)$,
$\omega_{1}(I)=J, \quad$ where $\quad J \in \sigma(\mathscr{D}),|J|=|I| / 4$ and $\inf J=\inf \sigma_{1}(I)$,
for all $I \in \mathscr{D}$. Now all we need to do is repeat the proof of Theorem 3.3 with $\omega$ replaced by $\omega_{\delta}$ in order to estimate $S_{\delta}$, for each $\delta \in\{0,1\}$.
4. The shift operator $T_{m}$. Here we define $16+4 \log _{2}(|m|), m \neq 0$, subcollections of the Haar system, so that on each such subcollection, $T_{m}$ acts as a martingale transform operator on either the dyadic grid or the one-thirdshifted dyadic grid. In Section 3 we established that changing the dyadic grid to the one-third-shifted dyadic grid is an isomorphism. Thus we may assume that $T_{m}$ is representable as a martingale transform operator on each of the $16+4 \log _{2}(|m|)$ subcollections, which yields the well known estimate

$$
\left\|T_{m}: L_{X}^{p} \rightarrow L_{X}^{p}\right\| \leq C\left(\log _{2}(2+|m|)\right)^{\alpha}
$$

for some $0<\alpha<1$, established by T. Figiel Fig88.
Define the shift map $\tau_{m}, m \in \mathbb{Z}$, by

$$
\begin{equation*}
\tau_{m}(I)=I+m|I| \tag{4.1}
\end{equation*}
$$

for all $I \in \mathscr{D} \cup \sigma(\mathscr{D})$. This induces the shift operator $T_{m}$, given by

$$
\begin{equation*}
T_{m} h_{I}=h_{\tau_{m}(I)} \tag{4.2}
\end{equation*}
$$

for all $I \in \mathscr{D} \cup \sigma(\mathscr{D})$. It is crucial that the one-third-shift operator $S$ defined in (3.4) and the shift operator $T_{m}$ commute, that is,

$$
\begin{equation*}
\left(S \circ T_{m}\right)(u)=\left(T_{m} \circ S\right)(u) \tag{4.3}
\end{equation*}
$$

for all $u \in L_{X}^{p}$. Analogously, we have

$$
\begin{align*}
& \left(S_{0} \circ T_{m}\right)(u)=\left(T_{m} \circ S_{0}\right)(u)  \tag{4.4}\\
& \left(S_{1} \circ T_{m}\right)(u)=\left(T_{m} \circ S_{1}\right)(u) \tag{4.5}
\end{align*}
$$

for all $u \in L_{X}^{p}($ see (3.9) $-(3.12))$.
We aim at splitting the dyadic intervals $\mathscr{D}$ into collections $\mathscr{B}_{i}^{(\delta)}$ such that we may bound $T_{m} \circ S^{\delta}$ on functions supported on $\sigma^{\delta}\left(\mathscr{B}_{i}^{(\delta)}\right), \delta \in\{0,1\}$. Note that if $\delta=0$, then $S^{\delta}=\mathrm{Id}$ and $\sigma^{\delta}=\mathrm{Id}$.

Given a shift width $m \in \mathbb{Z}, m \neq 0$, we will partition the dyadic intervals $\mathscr{D}$ into $16+4 \log _{2}(|m|)$ disjoint collections denoted by $\mathscr{B}_{i}^{(\delta)}$. The collections are constructed in such way that for each $i$ and $\delta \in\{0,1\}$ fixed, whenever $I \in \mathscr{B}_{i}^{(\delta)}$, the intervals $\sigma^{\delta}(I)$ and $\left(\tau_{m} \circ \sigma^{\delta}\right)(I)$ share the same dyadic predecessor with respect to the collection $\sigma^{\delta}\left(\mathscr{B}_{i}^{(\delta)}\right)$. The details are elaborated in Lemma 4.1 below.

Lemma 4.1. For every integer $m \in \mathbb{Z}, m \neq 0$ let $\tau_{m}$ denote the map given by

$$
\tau_{m}(I)=I+m|I|
$$

for all $I \in \mathscr{D} \cup \sigma(\mathscr{D})$. Then there exist a constant $K(m) \leq 7+2 \log _{2}(|m|)$ and disjoint collections of dyadic intervals $\mathscr{B}_{i}^{(\delta)}, 0 \leq i \leq K(m), \delta \in\{0,1\}$,
with

$$
\mathscr{D}=\bigcup_{\delta \in\{0,1\}} \bigcup_{i=0}^{K(m)} \mathscr{B}_{i}^{(\delta)},
$$

such that

$$
\begin{equation*}
\left\{I, \tau_{m}(I), I \cup \tau_{m}(I): I \in \sigma^{\delta}\left(\mathscr{B}_{i}^{(\delta)}\right)\right\} \tag{4.6}
\end{equation*}
$$

is a nested collection of sets, for all $0 \leq i \leq K(m)$ and $\delta \in\{0,1\}$.
Proof. By symmetry we may assume that $m \geq 1$, and we set $K(m)=$ $K(-m)$ if $m \leq-1$. So fix a shift width $m \geq 2$ and a $\lambda \geq 4$ such that

$$
\begin{equation*}
2^{\lambda-3} \leq m<2^{\lambda-2} \tag{4.7}
\end{equation*}
$$

and define $L(m)=\lambda-1$. If $m=1$, then let $\lambda=4$ and set $L(1)=3$. Now we split $\mathscr{D}$ into disjoint collections $\mathscr{A}_{i}, 0 \leq i \leq L(m)$, by omitting $L(m)$ consecutive levels of $\mathscr{D}$. More precisely, for every $0 \leq i \leq L(m)$ we define

$$
\begin{equation*}
\mathscr{A}_{i}=\bigcup_{j \in \mathbb{Z}}\left\{I \in \mathscr{D}:|I|=2^{-(\lambda j+i)}\right\} . \tag{4.8}
\end{equation*}
$$

Next we want to divide each $\mathscr{A}_{i}$ into two collections $\mathscr{A}_{i}^{(0)}$ and $\mathscr{A}_{i}^{(1)}$ so that every $I \in \mathscr{A}_{i}^{(0)}$ has the same predecessor in $\mathscr{A}_{i}^{(0)}$ as $\tau_{m}(I)$, and $\mathscr{A}_{i}^{(0)}$ is maximal. As a consequence, the collection $\mathscr{A}_{i}^{(1)}$ consists of all intervals $I$ such that $I$ and $\tau_{m}(I)$ do not share the same predecessor. But, if we apply the one-third-shift map $\sigma$ to $\mathscr{A}_{i}^{(1)}$, then every $I \in \sigma\left(\mathscr{A}_{i}^{(1)}\right)$ has the same predecessor in $\sigma\left(\mathscr{A}_{i}^{(1)}\right)$ as $\tau_{m}(I)$.

We will now construct these two collections. To this end let $\mathscr{G}$ denote one of the collections $\mathscr{A}_{i}, \sigma\left(\mathscr{A}_{i}\right), 0 \leq i \leq L(m)$, and define

$$
\begin{align*}
& \mathscr{C}_{0}(\mathscr{G}, I)=\left\{J \in \mathscr{G}:|J|=2^{-\lambda}|I|, J \subset I \text { and } \tau_{m}(J) \subset I\right\}, \\
& \mathscr{C}_{1}(\mathscr{G}, I)=\left\{J \in \mathscr{G}:|J|=2^{-\lambda}|I|, J \subset I \text { and } \tau_{m}(J) \cap I=\emptyset\right\} . \tag{4.9}
\end{align*}
$$

Revisiting the definition of the one-third-shift maps $\sigma, \sigma_{1}$ and considering the restriction (4.7) one can see that

$$
\begin{equation*}
\sigma\left(\mathscr{C}_{1}\left(\mathscr{A}_{i}, I\right)\right) \subset \mathscr{C}_{0}\left(\sigma\left(\mathscr{A}_{i}\right), \sigma_{1}(I)\right) \tag{4.10}
\end{equation*}
$$

for all $I \in \mathscr{A}_{i}, 0 \leq i \leq L(m)$. Note that for $I \in \mathscr{D}_{j}$ we have $\sigma_{1}(I)=\sigma(I)$ if $j$ is even and $\sigma_{1}(I)=\sigma(I+|I|)$ if $j$ is odd. Thus (4.10) implies that all intervals $J \in \sigma\left(\mathscr{C}_{1}\left(\mathscr{A}_{i}, I\right)\right)$ are such that $J$ and $\tau_{m}(J)$ share $\sigma_{1}(I)$ as common predecessor with respect to the collection $\sigma\left(\mathscr{A}_{i}^{(1)}\right)$. In Figure 1 one can see the action of the one-third-shift map $\sigma$ on the collection $\mathscr{A}_{i}$.


Fig. 1. The one-third-shift map $\sigma$ acting on $I \in \mathscr{D},|I|=2^{-j}$ and $J \in \mathscr{D},|J|=2^{-j-\lambda}$, where $J \subset I$ and $\tau_{m}(J) \cap I=\emptyset$. In this picture $\lambda$ is even.

Now define for every $0 \leq i \leq L(m)$ the following collections of dyadic intervals:

$$
\begin{align*}
& \mathscr{A}_{i}^{(0)}=\bigcup\left\{\mathscr{C}_{0}\left(\mathscr{A}_{i}, I\right): I \in \mathscr{A}_{i}\right\} \\
& \mathscr{A}_{i}^{(1)}=\mathscr{A}_{i} \backslash \mathscr{A}_{i}^{(0)} \tag{4.11}
\end{align*}
$$

Finally, for all $0 \leq i \leq L(m)$ and $\delta \in\{0,1\}$ we split $\mathscr{A}_{i}^{(\delta)}$ into two disjoint collections

$$
\begin{equation*}
\mathscr{B}_{i}^{(\delta)} \quad \text { and } \quad \mathscr{B}_{i+L(m)+1}^{(\delta)} \tag{4.12}
\end{equation*}
$$

such that

$$
\begin{equation*}
\mathscr{B}_{i}^{(\delta)} \cap \tau_{m}\left(\mathscr{B}_{i}^{(\delta)}\right)=\emptyset \tag{4.13}
\end{equation*}
$$

for all $0 \leq i \leq K(m)$ and $\delta \in\{0,1\}$, where we set $K(m)=2 L(m)+1$. Considering (4.7) and $L(m)=\lambda-1$ we find that $K(m) \leq 7+2 \log _{2}(m)$. For this purpose consider the collection

$$
\mathscr{E}=\left\{\tau_{k}(I): I \in \mathscr{D}, \inf I=0,0 \leq k \leq m-1\right\}
$$

and observe that

$$
\mathscr{D}=\bigcup_{\substack{j \in \mathbb{Z} \\ j \text { even }}} \tau_{j m}(\mathscr{E}) \cup \bigcup_{\substack{j \in \mathbb{Z} \\ j \text { odd }}} \tau_{j m}(\mathscr{E})=\mathscr{D}_{\text {even }} \cup \mathscr{D}_{\text {odd }}
$$

Now define the collections

$$
\begin{equation*}
\mathscr{B}_{i}^{(\delta)}=\mathscr{A}_{i}^{(\delta)} \cap \mathscr{D}_{\text {even }}, \quad \mathscr{B}_{i+L(m)+1}^{(\delta)}=\mathscr{A}_{i}^{(\delta)} \cap \mathscr{D}_{\text {odd }} \tag{4.14}
\end{equation*}
$$

for all $0 \leq i \leq L(m)$ and $\delta \in\{0,1\}$.
Taking into account (4.10), (4.9) and noting that $\tau_{m}(I) \in \mathscr{D}_{\text {odd }}$ if and only if $I \in \mathscr{D}_{\text {even }}$, we have verified (4.6), finishing this proof.

REmark 4.2. Note that we actually proved the slightly stronger result

$$
\begin{equation*}
I \cup \tau_{m}(I) \subset \pi^{\lambda}(I) \tag{4.15}
\end{equation*}
$$

for all $I \in \sigma^{\delta}\left(\mathscr{B}_{i}^{(\delta)}\right), 0 \leq i \leq K(m), \delta \in\{0,1\}$. For a given interval $I \in$ $\sigma^{\delta}(\mathscr{D}), \pi(I)$ is the unique interval $J \in \sigma^{\delta}(\mathscr{D})$ such that $J \supset I$ and $|J|=2|I|$.

As the combinatorial Lemma 4.1 exhibits the link between the shift map $\tau_{m}$, the one-third-shift map $\sigma$ and Figiel's compatibility condition (4.6), Theorem 4.3 below will translate the combinatorial results into analytical results, exhibiting the link between the shift operator $T_{m}$, the one-third-shift operator $S$ and martingale transform operators.

In the following context it is understood that $1<p<\infty, X$ is a Banach space with the UMD-property and $m \in \mathbb{Z}, m \neq 0$. Now we define projections $P_{i}^{(\delta)}: L_{X}^{p} \rightarrow L_{X}^{p}$, associated with the collections $\mathscr{B}_{i}^{(\delta)}$ of Lemma 4.1 by

$$
\begin{equation*}
P_{i}^{(\delta)} u=\sum_{I \in \mathscr{B}_{i}^{(\delta)}}\left\langle u, h_{I}\right\rangle h_{I}|I|^{-1} \tag{4.16}
\end{equation*}
$$

for all $0 \leq i \leq K(m), \delta \in\{0,1\}$ and $u \in L_{X}^{p}$. The UMD-property of $X$ implies uniform bounds on the projections $P_{i}^{(\delta)}$. Note the identity

$$
\begin{equation*}
u=\sum_{\delta \in\{0,1\}} \sum_{i=0}^{K(m)} P_{i}^{(\delta)} u \tag{4.17}
\end{equation*}
$$

for all $u \in L_{X}^{p}$, since the collections $\mathscr{B}_{i}^{(\delta)}, 0 \leq i \leq K(m), \delta \in\{0,1\}$, form a partition of $\mathscr{D}$ (see Lemma 4.1).

Exploiting the fact that the one-third-shift operator $S$ is an isomorphism on $L_{X}^{p}$ (see Theorem 3.2), we will now estimate the shift operator $T_{m}$ on the range of each $P_{i}^{(\delta)}$.

Theorem 4.3. Let $1<p<\infty$ and $X$ be a Banach space with the UMDproperty. Then for every $m \in \mathbb{Z}, 0 \leq i \leq K(m)$ and $\delta \in\{0,1\}$ the inequality

$$
\begin{equation*}
\left\|T_{m} \circ P_{i}^{(\delta)} u\right\|_{L_{X}^{p}} \leq C\left\|P_{i}^{(\delta)} u\right\|_{L_{X}^{p}} \tag{4.18}
\end{equation*}
$$

holds true for all $u \in L_{X}^{p}$, where the constant $C$ depends only on $\mathscr{U}_{p}(X)$. The projections $P_{i}^{(\delta)}, 0 \leq i \leq K(m), \delta \in\{0,1\}$, are defined according to (4.16), and $K(m) \leq 7+2 \log _{2}(1+|m|)$.

Proof. Note that by symmetry once we establish (4.18) for $m \geq 1$, the theorem will be proved.

Recalling the properties of the partition $\mathscr{B}_{i}^{(\delta)}, 0 \leq i \leq K(m), \delta \in\{0,1\}$, of $\mathscr{D}$ (see Lemma 4.1), we know that the collection

$$
\begin{equation*}
\left\{I, \tau_{m}(I), I \cup \tau_{m}(I): I \in \sigma^{\delta}\left(\mathscr{B}_{i}^{(\delta)}\right)\right\} \tag{4.19}
\end{equation*}
$$

is nested, for all $0 \leq i \leq K(m)$ and $\delta \in\{0,1\}$. Throughout this proof let $m \in \mathbb{Z}, 0 \leq i \leq K(m), \delta \in\{0,1\}$ and $u \in P_{i}^{(\delta)}\left(L_{X}^{p}\right)$ be fixed. According to (4.16) we may assume that $u$ has the representation

$$
u=\sum_{I \in \mathscr{B}_{i}^{(\delta)}} u_{I} h_{I}|I|^{-1} .
$$

For every $J \in \sigma^{\delta}(\mathscr{D})$ let

$$
\begin{equation*}
A^{(\delta)}(J)=J \cup \tau_{m}(J), \tag{4.20}
\end{equation*}
$$

and define

$$
\begin{equation*}
\mathscr{A}_{j}^{(\delta)}=\left\{A^{(\delta)}(J): J \in \sigma^{\delta}\left(\mathscr{D}_{j}\right) \cap \sigma^{\delta}\left(\mathscr{B}_{i}^{(\delta)}\right)\right\} \tag{4.21}
\end{equation*}
$$

for all $j \in \mathbb{Z}$. Then specify the filtration $\left\{\mathscr{F}_{j}^{(\delta)}\right\}_{j}$ by

$$
\begin{equation*}
\mathscr{F}_{j}^{(\delta)}=\sigma-\operatorname{alg}\left(\bigcup_{i \leq j} \mathscr{A}_{i}^{(\delta)}\right), \tag{4.22}
\end{equation*}
$$

and observe that due to (4.19) every $A^{(\delta)}(J), J \in \sigma^{\delta}\left(\mathscr{D}_{j}\right)$, is an atom for $\mathscr{F}_{j}^{(\delta)}$. The sigma algebra $\sigma-\operatorname{alg}\left(\bigcup_{i \leq j} \mathscr{A}_{i}^{(\delta)}\right)$ is the smallest sigma algebra containing $\bigcup_{i \leq j} \mathscr{A}_{i}^{(\delta)}$. The one-third-shift operator is given by

$$
\begin{equation*}
S^{\delta} u=\sum_{I \in \mathscr{B}_{i}^{(\delta)}} u_{I} h_{\sigma^{\delta}(I)}|I|^{-1}=\sum_{J \in \sigma^{\delta}\left(\mathscr{B}_{i}^{(\delta)}\right)} u_{\sigma^{-\delta}(J)} h_{J}|J|^{-1} \tag{4.23}
\end{equation*}
$$

(see (3.4) for details). We recall the notation

$$
(u)_{j}=\sum_{|I|=2^{-j}} u_{I} h_{I}|I|^{-1} \quad \text { and } \quad \mathbb{I}(u)_{j}=\sum_{|I|=2^{-j}} u_{I} 1_{I}|I|^{-1},
$$

and note that

$$
\begin{equation*}
\left\|T_{m} S^{\delta} u\right\|_{L_{X}^{p}} \approx \int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t) \mathbb{I}\left(T_{m} S^{\delta} u\right)_{j}\right\|_{L_{X}^{p}} d t \tag{4.24}
\end{equation*}
$$

(see (3.6)-(3.8)). If $I \in \mathscr{D}_{j} \cap \mathscr{B}_{i}^{(\delta)}$ then

$$
\mathbb{E}\left(1_{\sigma^{\delta}(I)} \mid \mathscr{F}_{j}^{(\delta)}\right)=\left|A^{(\delta)}\left(\sigma^{\delta}(I)\right)\right|^{-1} 1_{A^{(\delta)}\left(\sigma^{\delta}(I)\right)} \geq \frac{1}{2} 1_{\tau_{m}\left(\sigma^{\delta}(I)\right)},
$$

therefore

$$
\mathbb{I}\left(T_{m} S^{\delta} u\right)_{j} \leq 2 \mathbb{E}\left(\mathbb{I}\left(S^{\delta} u\right)_{j} \mid \mathscr{F}_{j}^{(\delta)}\right)
$$

Thus, Kahane's contraction principle and Bourgain's version of Stein's martingale inequality yield

$$
\begin{aligned}
\int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t) \mathbb{I}\left(T_{m} S^{\delta} u\right)_{j}\right\|_{L_{X}^{p}} d t & \leq \int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t) 2 \mathbb{E}\left(\mathbb{I}\left(S^{\delta} u\right)_{j} \mid \mathscr{F}_{j}^{(\delta)}\right)\right\|_{L_{X}^{p}} d t \\
& \lesssim \int_{0}^{1}\left\|\sum_{j \in \mathbb{Z}} r_{j}(t) \mathbb{I}\left(S^{\delta} u\right)_{j}\right\|_{L_{X}^{p}} d t \\
& \approx\left\|S^{\delta} u\right\|_{L_{X}^{p}} .
\end{aligned}
$$

Combining the latter estimate with (4.24), Theorem 3.2 proves

$$
\begin{equation*}
\left\|T_{m} S^{\delta} u\right\|_{L_{X}^{p}} \lesssim\|u\|_{L_{X}^{p}} \tag{4.25}
\end{equation*}
$$

According to (4.3) the shift operator $T_{m}$ and the one-third-shift operator $S$ commute, so we have the identity

$$
T_{m} u=\left(S^{-\delta} \circ T_{m} \circ S^{\delta}\right)(u),
$$

and we obtain, by an application of Theorem 3.2,

$$
\begin{equation*}
\left\|T_{m} u\right\|_{L_{X}^{p}} \lesssim\left\|\left(T_{m} \circ S^{\delta}\right)(u)\right\|_{L_{X}^{p}} . \tag{4.26}
\end{equation*}
$$

We conclude the proof by combining (4.26) and (4.25).
REMARK 4.4. By slightly adjusting the construction of $\mathscr{B}_{i}^{(\delta)}$ we could replace Bourgain's version of Stein's martingale inequality by the martingale transforms in (Fig88, Proposition 2, Step 0] in order to obtain (4.25).

Within this remark we shall make use of the following notation. If $I \in$ $\mathscr{D} \cup \sigma(\mathscr{D})$, then $I_{0}$ is the uniquely defined interval $J \in \mathscr{D} \cup \sigma(\mathscr{D})$ with $|J|=|I| / 2$ such that $\inf J=\inf I$, and $I_{1}$ is the uniquely defined interval $J \in \mathscr{D} \cup \sigma(\mathscr{D})$ with $|J|=|I| / 2$ such that $\sup J=\sup I$.

Let us now describe the aforementioned modifications. We have to replace $\lambda$ by $\lambda+1$ and redefine $\mathscr{C}_{0}$ and $\mathscr{C}_{1}$ as

$$
\begin{aligned}
\mathscr{C}_{0}\left(I, \mathscr{A}_{i}\right)= & \left\{J \in \mathscr{A}_{i}:|J|=2^{-\lambda}|I|, J \subset I_{0} \text { and } \tau_{m}(J) \subset I_{0}\right\} \\
& \cup\left\{J \in \mathscr{A}_{i}:|J|=2^{-\lambda}|I|, J \subset I_{1} \text { and } \tau_{m}(J) \subset I_{1}\right\}, \\
\mathscr{C}_{1}\left(I, \mathscr{A}_{i}\right)= & \left\{J \in \mathscr{A}_{i}:|J|=2^{-\lambda}|I|, J \subset I_{0} \text { and } \tau_{m}(J) \cap I_{0}=\emptyset\right\} \\
& \cup\left\{J \in \mathscr{A}_{i}:|J|=2^{-\lambda}|I|, J \subset I_{1} \text { and } \tau_{m}(J) \cap I_{1}=\emptyset\right\}
\end{aligned}
$$

(cf. (4.8) and (4.9)). This results in the collection

$$
\begin{equation*}
\left\{J_{0}, \tau_{m}(J)_{0}, J_{1}, \tau_{m}(J)_{1}, J \cup \tau_{m}(J): J \in \sigma^{\delta}\left(\mathscr{B}_{i}^{(\delta)}\right)\right\} \tag{4.27}
\end{equation*}
$$

being nested for all $0 \leq i \leq K(m)$ and $\delta \in\{0,1\}$. With these modifications let us define

$$
d_{J, 1}^{(\delta)}=\frac{1}{2}\left(h_{J}+h_{\tau_{m}(J)}\right) \quad \text { and } \quad d_{J, 2}^{(\delta)}=\frac{1}{2}\left(h_{J}-h_{\tau_{m}(J)}\right),
$$

for all $J \in \sigma^{\delta}\left(\mathscr{B}_{i}^{(\delta)}\right)$. Since (4.27) is nested, $\left\{d_{J, 1}^{(\delta)}, d_{J, 2}^{(\delta)}: J \in \sigma\left(\mathscr{B}_{i}^{(\delta)}\right)\right\}$ forms a martingale difference sequence. Observe $h_{J}=d_{J, 1}^{(\delta)}+d_{J, 2}^{(\delta)}$ and $h_{\tau_{m}(J)}=$ $d_{J, 1}^{(\delta)}-d_{J, 2}^{(\delta)}$, hence we may swap $h_{J}$ and $h_{\tau_{m}(J)}$ without using Bourgain's version of Stein's martingale inequality.
5. A martingale decomposition for $U_{m}$. Recall that in (4.1) we defined the shift map $\tau_{m}$ for every $m \in \mathbb{Z}$ by

$$
\tau_{m}(I)=I+m|I|
$$

for all $I \in \mathscr{D} \cup \sigma(\mathscr{D})$. We introduce the operator $U_{m}$ by setting

$$
\begin{equation*}
U_{m} h_{I}=1_{\tau_{m}(I)}-1_{I} \tag{5.1}
\end{equation*}
$$

for all $I \in \mathscr{D} \cup \sigma(\mathscr{D})$. Although the operators $T_{m}$ and $U_{m}$ appear to be similar, $U_{m}$ is much harder to analyze than $T_{m}$. This is rooted in the observation that $\left\{T_{m} h_{I}\right\}_{I \in \mathscr{A}}$ is a martingale difference sequence for any choice of $\mathscr{A} \subset \mathscr{D}$, whereas whether $\left\{U_{m} h_{I}\right\}_{I \in \mathscr{B}}$ forms a martingale difference sequence strongly depends on the choice of $\mathscr{B} \subset \mathscr{D}$. Essentially the same method we used to bound $T_{m}$ for functions supported on the collections $\mathscr{B}_{i}^{(0)}, 0 \leq i \leq K(m)$, qualifies for estimating $U_{m}$, since Lemma 4.1 ensures that $\left\{U_{m} h_{I}: I \in \mathscr{B}_{i}^{(0)}\right\}$ forms a martingale difference sequence. So the main obstacle is to estimate $U_{m}$ on $\mathscr{B}_{i}^{(1)}$, since $\left\{U_{m} h_{I}: I \in \mathscr{B}_{i}^{(1)}\right\}$ is not a martingale difference sequence.

The remedy to this problem is the martingale difference sequence decomposition of $U_{m}$ into

$$
U_{m}=a_{I}+b_{I}+c_{I}, \quad I \in \mathscr{B}_{i}^{(1, \varepsilon)},
$$

for suitable subcollections $\mathscr{B}_{i}^{(1, \varepsilon)}, \varepsilon \in\{0,1\}$, of $\mathscr{B}_{i}^{(1)}$, such that the collections $\left\{a_{I}: I \in \mathscr{B}_{i}^{(1, \varepsilon)}\right\},\left\{b_{I}: I \in \mathscr{B}_{i}^{(1, \varepsilon)}\right\}$ and $\left\{c_{I}: I \in \mathscr{B}_{i}^{(1, \varepsilon)}\right\}$ form martingale difference sequences. Since there are $56+14 \log _{2}(|m|)$ subcollections $\mathscr{B}_{i}^{(\delta, \varepsilon)}$, we can obtain the estimate

$$
\left\|U_{m}: L_{X}^{p} \rightarrow L_{X}^{p}\right\| \leq C\left(\log _{2}(2+|m|)\right)^{\beta} \quad \text { for some } 0<\beta<1,
$$

established in Fig88 by the same method we used for $T_{m}$ in Section [4.
First, let us define the maps $\beta_{0}, \beta_{1}$ and $\beta$ by

$$
\begin{align*}
\beta_{0}(I) & =\sigma_{0}(I) \backslash I,  \tag{5.2}\\
\beta_{1}(I) & =\sigma_{1}(I) \cap I,  \tag{5.3}\\
\beta(I) & =\beta_{0}(I) \cup \beta_{1}(I), \tag{5.4}
\end{align*}
$$

where $\sigma_{0}$ and $\sigma_{1}$ are given by (3.9) and (3.10). Secondly, set

$$
\begin{align*}
\gamma_{0}(I) & =\sigma_{0}(I) \cap I  \tag{5.5}\\
\gamma_{1}(I) & =\sigma_{1}(I) \backslash I  \tag{5.6}\\
\gamma(I) & =\gamma_{0}(I) \cup \gamma_{1}(I) . \tag{5.7}
\end{align*}
$$

The functions $\sigma_{0}, \sigma_{1}, \beta_{0}, \beta_{1}$ and $\gamma_{0}, \gamma_{1}$ are visualized in Figure 2,


Fig. 2. The functions $\sigma_{0}, \sigma_{1}, \beta_{0}, \beta_{1}$ and $\gamma_{0}, \gamma_{1}$ applied to $I$ or $\tau_{m}(I)$

Similar to the combinatorial Lemma 4.1 for the operator $T_{m}$ we shall now display a combinatorial lemma for $U_{m}$, which will enable us to use the same method to estimate $U_{m}$ as well. The estimates of $U_{m}$ will be established in Theorem 5.2,

LEMmA 5.1. Let $m \in \mathbb{Z}, m \neq 0$ and let $\mathscr{B}_{i}^{(\delta)}, 0 \leq i \leq K(m), \delta \in\{0,1\}$, and $K(m)$ be defined as in Lemma4.1. If we split each collection $\mathscr{B}_{i}^{(1)}$ into

$$
\begin{equation*}
\mathscr{B}_{i}^{(1, \varepsilon)}=\mathscr{B}_{i}^{(1)} \cap\left\{\tau_{k}(I): I \in \mathscr{D}, \inf I=0, k \bmod 2=\varepsilon\right\} \tag{5.8}
\end{equation*}
$$

for all $0 \leq i \leq K(m), \varepsilon \in\{0,1\}$, then

$$
\begin{align*}
& \left\{I, \tau_{m}(I), I \cup \tau_{m}(I): I \in \mathscr{B}_{i}^{(0)}\right\}  \tag{5.9}\\
& \left\{\sigma_{0}(I), \sigma_{1}\left(\tau_{m}(I)\right), \sigma_{0}(I) \cup \sigma_{1}\left(\tau_{m}(I)\right): I \in \mathscr{B}_{i}^{(1, \varepsilon)}\right\},  \tag{5.10}\\
& \left\{\beta_{0}(I), \beta_{1}(I), \beta_{0}(I) \cup \beta_{1}(I): I \in \mathscr{B}_{i}^{(1, \varepsilon)}\right\},  \tag{5.11}\\
& \left\{\gamma_{0}(I), \gamma_{1}(I), \gamma_{0}(I) \cup \gamma_{1}(I): I \in \tau_{m}\left(\mathscr{B}_{i}^{(1, \varepsilon)}\right)\right\} \tag{5.12}
\end{align*}
$$

are nested collections of sets, for all $0 \leq i \leq K(m)$ and $\varepsilon \in\{0,1\}$.
Proof. By symmetry we may assume $m \geq 1$. Let $m \in \mathbb{Z}, m \geq 1,0 \leq$ $i \leq K(m)$ and $\varepsilon \in\{0,1\}$ be fixed throughout the rest of this proof. Given $I \in \sigma^{\delta}(\mathscr{D})$ its predecessor $\pi(I)$ is the uniquely determined interval $J \in \sigma^{\delta}(\mathscr{D})$ such that $J \supset I$ and $|J|=2|I|$. Given intervals $J \in \mathscr{B}_{i}^{(1)}$ and $K=\pi^{\lambda}(J)$
observe that

$$
\begin{equation*}
\tau_{m}(J) \cap K=\emptyset, \quad \inf J \geq \inf K+\frac{3}{4}|K|, \quad \tau_{-1}(J) \subset K \tag{5.13}
\end{equation*}
$$

First, note that (5.9) is covered by Lemma 4.1.
Secondly, we will show that (5.10) is nested. Henceforth, we shall abbreviate $\mathscr{B}_{i}^{(1, \varepsilon)}$ by $\mathscr{B}$. Let $I, J \in \mathscr{B}$ be such that

$$
\left(\sigma_{0}(J) \cup \sigma_{1}\left(\tau_{m}(J)\right)\right) \cap\left(\sigma_{0}(I) \cup \sigma_{1}\left(\tau_{m}(I)\right)\right) \neq \emptyset
$$

Note that if $|I|=|J|$, then $I=J$ by definition of $\mathscr{B}$. So assume that $|J|<|I|$. Defining $K:=\pi^{\lambda}(J)$ and recalling (5.13) we see $\inf J \geq \inf K+\frac{3}{4}|K|$, therefore $\sigma_{0}(J) \subset \sigma_{1}(K)$. Furthermore, considering the constraints for the parameters $\lambda$ and $m$ we obtain

$$
\sup \sigma_{1}\left(\tau_{m}(J)\right) \leq \sup K+\frac{2}{3}|J|+m|J|<\sup K+\frac{1}{3}|K|
$$

thus $\sigma_{1}\left(\tau_{m}(J)\right) \subset \sigma_{1}(K)$ as well. So we have proved that $\sigma_{0}(J) \cup \sigma_{1}(J) \subset$ $\sigma_{1}(K)$. The nestedness of the one-third-shifted intervals implies either $\sigma_{1}(K)$ $\subset \sigma_{0}(I)$ or $\sigma_{1}(K) \subset \sigma_{1}(I)$, hence

$$
\sigma_{0}(J) \cup \sigma_{1}(J) \subset \sigma_{0}(I) \quad \text { or } \quad \sigma_{0}(J) \cup \sigma_{1}(J) \subset \sigma_{1}(I)
$$

Thirdly, we prove the nestedness of (5.11). To this end, let $I, J \in \mathscr{B}$ be such that

$$
\beta(J) \cap \beta(I) \neq \emptyset .
$$

As in the proof of (5.10) we see that $|I|=|J|$ implies $I=J$. Let us now assume $|J|<|I|$,

$$
\begin{equation*}
\beta(J) \cap \beta(I) \neq \emptyset \quad \text { and } \quad \beta(J) \cap \beta(I)^{c} \neq \emptyset \tag{A}
\end{equation*}
$$

Considering that $\beta(I)=\beta_{0}(I) \cup \beta_{1}(I)$ and $\beta(J) \subset J \cup \tau_{-1}(J)$, assumption (A) is covered by the following four cases:
(1) $\inf \beta_{0}(I) \in J \cup \tau_{-1}(J)$,
(2) $\sup \beta_{0}(I) \in J \cup \tau_{-1}(J)$,
(3) $\inf \beta_{1}(I) \in J \cup \tau_{-1}(J)$,
(4) $\sup \beta_{1}(I) \in J \cup \tau_{-1}(J)$.

Assume case (1) is true. Define $K=\pi^{\lambda}(J)$ and note that $|K| \leq|I|$ by definition of $\mathscr{B}$. Since the endpoint of the one-third-shifted interval $\sigma_{0}(I)$ is contained in a non-shifted smaller interval $K$, we obtain

$$
\begin{equation*}
\inf \beta_{0}(I)=\inf K+\frac{z}{3}|K| \tag{5.14}
\end{equation*}
$$

for some $z \in\{1,2\}$, depending on the direction of the one-third-shift. We know from (5.13) that $\inf J \geq \inf K+\frac{3}{4}|K|$, hence

$$
\inf J \cup \tau_{-1}(J) \geq \inf K+\frac{11}{16}|K|>\inf K+\frac{2}{3}|K|
$$

which contradicts (5.14). If we assume case (2) is true, then $\sup \beta_{0}(I)=$ $\inf I$ and we find that $\inf J=\inf I$. Consequently, $\inf J=\inf \pi^{\lambda}(J)$, which contradicts $J \in \mathscr{B}$. Case (3) is analogous to (11), and case (4) is analogous to (2). Altogether we have proved that assumption ( $(\mathcal{A})$ is false, therefore, considering that $\operatorname{dist}\left(\beta_{0}(I), \beta_{1}(I)\right) \geq \frac{1}{3}|I|$ and $\operatorname{diam}(\beta(J)) \leq \frac{1}{8}|I|$ yields

$$
\beta(J) \subset \beta_{0}(I) \quad \text { or } \quad \beta(J) \subset \beta_{1}(I)
$$

for all $I, J \in \mathscr{B}$ with $|J|<|I|$ such that $\beta(J) \cap \beta(I) \neq \emptyset$.
The proof that (5.12) is nested uses essentially the same argument presented in the proof of (5.11), therefore we omit the details. Both of these cases use that if $J \in \tau_{m}(\mathscr{B})$, then $\gamma(J) \subset J \cup \tau_{1}(J), \pi^{\lambda}(J)=\pi^{\lambda}\left(\tau_{1}(J)\right)$ and $\inf \tau_{1}(J) \neq \inf \pi^{\lambda}\left(\tau_{1}(J)\right)$.

With $m \in \mathbb{Z}$ fixed, we introduce the functions

$$
\begin{array}{ll}
a_{I}=1_{\sigma_{1}\left(\tau_{m}(I)\right)}-1_{\sigma_{0}(I)}, & \\
b_{I}=1_{\beta_{0}(I)}-1_{\beta_{1}(I)}, & \\
c_{I}=1_{\gamma_{0}\left(\tau_{m}(I)\right)}-1_{\gamma_{1}\left(\tau_{m}(I)\right)}, &  \tag{5.17}\\
I \in \mathscr{D}, \\
\end{array}
$$

(see Figure 3).


Fig. 3. Martingale decomposition of $U_{m}$

Accordingly, we define operators $\mathcal{A}, \mathcal{B}$ and $\mathcal{C}$ as the linear extensions of

$$
\begin{array}{ll}
\mathcal{A} h_{I}=a_{I}, & I \in \mathscr{D}, \\
\mathcal{B} h_{I}=b_{I}, & I \in \mathscr{D}, \\
\mathcal{C} h_{I}=c_{I}, & I \in \mathscr{D} . \tag{5.20}
\end{array}
$$

Note the identities

$$
\begin{align*}
U_{m} h_{I} & =a_{I}+b_{I}+c_{I}, & & I \in \mathscr{D}, \\
U_{m} u & =\mathcal{A} u+\mathcal{B} u+\mathcal{C} u, & & u \in L_{X}^{p},  \tag{5.21}\\
\mathcal{A} u & =\left(U_{m+1} \circ S_{0}\right)(u), & & u \in L_{X}^{p},
\end{align*}
$$

as well as the estimates

$$
\begin{array}{ll}
\left|b_{I}\right| \leq\left|h_{I}\right|+\left|S_{0}\left(h_{I}\right)\right|, & I \in \mathscr{D}, \\
\left|c_{I}\right| \leq\left|T_{m}\left(S_{0} h_{I}\right)\right|+\left|T_{m}\left(S_{1} h_{I}\right)\right|, & I \in \mathscr{D} \tag{5.22}
\end{array}
$$

(see Figure (3).
We will introduce projections $P_{i}^{(0)}, P_{i}^{(1, \varepsilon)}, 0 \leq i \leq K(m), \varepsilon \in\{0,1\}$, and estimate the operators $\mathcal{A}, \mathcal{B}$ and $\mathcal{C}$ on their respective images (see Theorem 5.22). Recall that

$$
P_{i}^{(0)} u=\sum_{I \in \mathscr{B}_{i}^{(0)}}\left\langle u, h_{I}\right\rangle h_{I}|I|^{-1}, \quad 0 \leq i \leq K(m),
$$

was already defined in (4.16). Accordingly, we set

$$
\begin{equation*}
P_{i}^{(1, \varepsilon)} u=\sum_{I \in \mathscr{B}_{i}^{(1, \varepsilon)}}\left\langle u, h_{I}\right\rangle h_{I}|I|^{-1}, \quad 0 \leq i \leq K(m), \varepsilon \in\{0,1\}, \tag{5.23}
\end{equation*}
$$

for all $u \in L_{X}^{p}$, where $\mathscr{B}_{i}^{(1, \varepsilon)}$ is defined in (5.8). Note that $P_{i}^{(1)}=P_{i}^{(1,0)}+$ $P_{i}^{(1,1)}$, where $P_{i}^{(1)}$ was defined in (4.16). Thus, (4.17) implies the identity

$$
\begin{equation*}
u=\sum_{i=0}^{K(m)} P_{i}^{(0)} u+P_{i}^{(1,0)} u+P_{i}^{(1,1)} u, \quad u \in L_{X}^{p} . \tag{5.24}
\end{equation*}
$$

Consider the splitting of $\mathscr{D}$ into the sets $\mathscr{B}_{i}^{(\delta)}, 0 \leq i \leq K(m), \delta \in\{0,1\}$ (see Lemma 4.1 for details), which we used in Theorem 4.3 to treat the operator $T_{m}$. Retracing our steps in the proof of Theorem 4.3 we shall find that with the result in Lemma 5.1 we could actually repeat this proof with the operator $T_{m} \circ P_{i}^{(\delta)}, \delta \in\{0,1\}$, replaced by one of the operators $U_{m} \circ P_{i}^{(0)}$ or $\mathcal{A} \circ P_{i}^{(1, \varepsilon)}, \varepsilon \in\{0,1\}$. The residual terms $\mathcal{B} \circ P_{i}^{(1, \varepsilon)}$ and $\mathcal{C} \circ P_{i}^{(1, \varepsilon)}$ occurring in the martingale decomposition of $U_{m}$ are dominated by $\left(\operatorname{Id}+S_{0}\right) \circ P_{i}^{(1, \varepsilon)}$ and $\left(T_{m} \circ\left(S_{0}+S_{1}\right)\right) \circ P_{i}^{(1, \varepsilon)}$, respectively.

Theorem 5.2. Let $m \in \mathbb{Z}, m \neq 0$ and let $P_{i}^{(0)}, P_{i}^{(1, \varepsilon)}, \varepsilon \in\{0,1\}$, be defined according to (4.16) and (5.23). Then

$$
\begin{equation*}
U_{m} u=\sum_{i=0}^{K(m)}\left\{\left(U_{m} \circ P_{i}^{(0)}\right)(u)+\sum_{\varepsilon \in\{0,1\}}\left((\mathcal{A}+\mathcal{B}+\mathcal{C}) \circ P_{i}^{(1, \varepsilon)}\right)(u)\right\} \tag{5.25}
\end{equation*}
$$

for all $u \in L_{X}^{p}$. Furthermore, we have the estimates

$$
\begin{align*}
\left\|\left(U_{m} \circ P_{i}^{(0)}\right)(u)\right\|_{L_{X}^{p}} & \leq C\left\|P_{i}^{(0)} u\right\|_{L_{X}^{p}}  \tag{5.26}\\
\left\|\left(\mathcal{A} \circ P_{i}^{(1, \varepsilon)}\right)(u)\right\|_{L_{X}^{p}} & \leq C\left\|P_{i}^{(1, \varepsilon)} u\right\|_{L_{X}^{p}}  \tag{5.27}\\
\left\|\left(\mathcal{B} \circ P_{i}^{(1, \varepsilon)}\right)(u)\right\|_{L_{X}^{p}} & \leq C\left\|P_{i}^{(1, \varepsilon)} u\right\|_{L_{X}^{p}}  \tag{5.28}\\
\left\|\left(\mathcal{C} \circ P_{i}^{(1, \varepsilon)}\right)(u)\right\|_{L_{X}^{p}} & \leq C\left\|P_{i}^{(1, \varepsilon)} u\right\|_{L_{X}^{p}} \tag{5.29}
\end{align*}
$$

for all $u \in L_{X}^{p}, 0 \leq i \leq K(m)$ and $\varepsilon \in\{0,1\}$. The constant $C$ depends only on $\mathscr{U}_{p}(X)$ and we have the estimate $K(m) \leq 7+2 \log _{2}(|m|)$.

Proof. Let $m \in \mathbb{Z}, m \neq 0,0 \leq i \leq K(m)$ and $\varepsilon \in\{0,1\}$ be fixed throughout the proof. By symmetry we may assume that $m \geq 1$.

Note that (5.21) and (5.24) imply identity (5.25).
From Lemma 5.1 we know that each of the collections

$$
\begin{aligned}
\left\{U_{m} h_{I}: I \in \mathscr{B}_{i}^{(0)}\right\}, & \left\{\mathcal{A} h_{I}: I \in \mathscr{B}_{i}^{(1, \varepsilon)}\right\}, \\
\left\{\mathcal{B} h_{I}: I \in \mathscr{B}_{i}^{(1, \varepsilon)}\right\}, & \left\{\mathcal{C} h_{I}: I \in \mathscr{B}_{i}^{(1, \varepsilon)}\right\}
\end{aligned}
$$

forms a martingale difference sequence. Hence, in order to estimate $\left\{U_{m} h_{I}\right.$ : $\left.I \in \mathscr{B}_{i}^{(0)}\right\}$ and $\left\{\mathcal{A} h_{I}: I \in \mathscr{B}_{i}^{(1, \varepsilon)}\right\}$ we can essentially repeat the proof of Theorem4.3, thereby obtaining (5.26) and (5.27). Since $\left|b_{I}\right| \leq\left|h_{I}\right|+\left|S_{0}\left(h_{I}\right)\right|$ (see (5.22)), we deduce from the UMD-property, Kahane's inequality (2.3) and Kahane's contraction principle (2.2) that

$$
\left\|\left(\mathcal{B} \circ P_{i}^{(1, \varepsilon)}\right)(u)\right\|_{L_{X}^{p}} \lesssim\left\|P_{i}^{(1, \varepsilon)} u\right\|_{L_{X}^{p}}+\left\|\left(S_{0} \circ P_{i}^{(1, \varepsilon)}\right)(u)\right\|_{L_{X}^{p}} .
$$

The latter estimate together with Theorem 3.3 implies (5.28). Observe $\left|c_{I}\right| \leq$ $\left|T_{m}\left(S_{0} h_{I}\right)\right|+\left|T_{m}\left(S_{1} h_{I}\right)\right|$ (see (5.22)), thus the UMD-property, Kahane's inequality and Kahane's contraction principle imply
$\left\|\left(\mathcal{C} \circ P_{i}^{(1, \varepsilon)}\right)(u)\right\|_{L_{X}^{p}} \lesssim\left\|\left(S_{0} \circ T_{m} \circ P_{i}^{(1, \varepsilon)}\right)(u)\right\|_{L_{X}^{p}}+\left\|\left(S_{1} \circ T_{m} \circ P_{i}^{(1, \varepsilon)}\right)(u)\right\|_{L_{X}^{p}}$.
Note that (5.10) implies that

$$
\left\{\sigma_{\delta}(I), \sigma_{\delta}\left(\tau_{m}(I)\right), \sigma_{\delta}(I) \cup \sigma_{\delta}\left(\tau_{m}(I)\right): I \in \mathscr{B}_{i}^{(1, \varepsilon)}\right\}
$$

is nested, hence we can replicate the proof of Theorem 4.3 once more.
From Theorems 4.3 and 5.2, one can obtain the estimates stated in Theorem 5.3 below by exploiting the type and cotype inequalities for $T_{m}$, and only the cotype inequality for $U_{m}$. Inserting the estimates of Theorems 4.3
and 5.2 into Fig88, Lemma 1] one can obtain Fig88, Theorem 1], stated below for the sake of completeness.

Theorem 5.3 (Fig88). Let $1<p<\infty$, and $X$ be a Banach space with the UMD-property. For $m \in \mathbb{Z}$ let $\tau_{m}$ denote the shift map defined by

$$
I \mapsto I+m|I|
$$

Let $T_{m}, U_{m}$ denote the linear extensions of the maps

$$
T_{m} h_{I}=h_{\tau_{m}(I)}
$$

and

$$
U_{m} h_{I}=1_{\tau_{m}(I)}-1_{I},
$$

respectively. Then

$$
\begin{align*}
\left\|T_{m}: L_{X}^{p} \rightarrow L_{X}^{p}\right\| \leq C\left(\log _{2}(2+|m|)\right)^{\alpha}  \tag{5.30}\\
\left\|U_{m}: L_{X}^{p} \rightarrow L_{X}^{p}\right\| \leq C\left(\log _{2}(2+|m|)\right)^{\beta} \tag{5.31}
\end{align*}
$$

where the constant $C>0$ depends only on $\mathscr{U}_{p}(X)$ and $0<\alpha, \beta<1$. Moreover, if $L_{X}^{p}$ has type $\mathcal{T}$ and cotype $\mathfrak{C}$, then one can take $\alpha=1 / \mathcal{T}-1 / \mathcal{C}$ and $\beta=1-1 / \mathcal{C}$.

These bounds are used to establish equivalence of some bases to the Haar system, e.g. the Franklin system (cf. Fig88). In Fig90, T. Figiel bounds singular integral operators $S$ acting on vector-valued $L^{p}$ spaces by means of the estimates (5.30) and (5.31). This is achieved by decomposing $S$ into weighted series of the rearrangement operators $T_{m}$ and $U_{m}$.

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