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## Full Length Article

# When are the norms of the Riesz projection and the backward shift operator equal to one? 

Oleksiy Karlovych ${ }^{\text {a,* }}$, Eugene Shargorodsky ${ }^{\text {b }}$<br>${ }^{\text {a }}$ Centro de Matemática e Aplicações, Departamento de Matematica, Faculdade de Ciencias e Tecnologia, Universidade Nova de Lisboa, Quinta da Torre, Caparica, 2829-516, Portugal<br>${ }^{\mathrm{b}}$ Department of Mathematics, King's College London, Strand, London, WC2R 2LS, United Kingdom

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The lower estimate by Gohberg and Krupnik (1968) and the upper estimate by Hollenbeck and Verbitsky (2000) for the norm of the Riesz projection $P$ on the Lebesgue space $L^{p}$ lead to $\|P\|_{L^{p} \rightarrow L^{p}}=1 / \sin (\pi / p)$ for every $p \in(1, \infty)$. Hence $L^{2}$ is the only space among all Lebesgue spaces $L^{p}$ for which the norm of the Riesz projection $P$ is equal to one. Banach function spaces $X$ are far-reaching generalisations of Lebesgue spaces $L^{p}$. We prove that the norm of $P$ is equal to one on the space $X$ if and only if $X$ coincides with $L^{2}$ and there exists a constant $C \in(0, \infty)$ such that $\|f\|_{X}=C\|f\|_{L^{2}}$ for all functions $f \in X$. Independently from this, we also show that the norm of $P$ on $X$ is equal to one if and only if the norm of the backward shift operator $S$ on the abstract Hardy space $H[X]$ built upon $X$ is equal to one.
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## 1. Introduction

For a function $f \in L^{1}$ on the unit circle $\mathbb{T}:=\{z \in \mathbb{C}:|z|=1\}$, let

$$
\widehat{f}(n)=\frac{1}{2 \pi} \int_{-\pi}^{\pi} f\left(e^{i \theta}\right) e^{-i n \theta} d \theta, \quad n \in \mathbb{Z}
$$

be the Fourier coefficients of $f$. Let $X$ be a Banach function space on $\mathbb{T}$. We postpone the technical definition until Section 2.1 and mention here only that $X$ is continuously embedded into $L^{1}$. Let

$$
H[X]:=\{g \in X: \widehat{g}(n)=0 \quad \text { for all } \quad n<0\}
$$

denote the abstract Hardy space built upon the Banach function space $X$. In the case $X=L^{p}$, we will use the standard notation $H^{p}:=H\left[L^{p}\right]$. We will also use the following notation:

$$
\mathbf{e}_{m}(z):=z^{m}, \quad z \in \mathbb{C}, \quad m \in \mathbb{Z}
$$

It is easy to see that the backward shift operator $S$, defined by

$$
(S f)(t):=\mathbf{e}_{-1}(t)(f(t)-\widehat{f}(0)), \quad t \in \mathbb{T}
$$

is bounded on the space $H[X]$. Consider the operators $\mathcal{C}$ and $P$, defined for a function $f \in L^{1}$ and an a.e. point $t \in \mathbb{T}$ by

$$
(\mathcal{C} f)(t):=\frac{1}{\pi i} \text { p.v. } \int_{\mathbb{T}} \frac{f(\tau)}{\tau-t} d \tau, \quad(P f)(t):=\frac{f(t)+(\mathcal{C} f)(t)}{2}
$$

respectively, where the integral is understood in the Cauchy principal value sense. The operator $\mathcal{C}$ is called the Cauchy singular integral operator and the operator $P$ is called the Riesz projection. The latter term can be explained by the fact that if $P$ is bounded on a Banach function space $X$, then one has $H[X]=P(X)$ (see [23, Lemma 1.1]).

The lower estimate by Gohberg and Krupnik (see [16, Ch. 9, Theorem 9.1]) and the upper estimate by Hollenbeck and Verbitsky [18] for the norm of the Riesz projection $P$ on the Lebesgue space $L^{p}$ lead to

$$
\begin{equation*}
\|P\|_{L^{p} \rightarrow L^{p}}=1 / \sin (\pi / p), \quad 1<p<\infty . \tag{1.1}
\end{equation*}
$$

Hence $L^{2}$ is the only space among all Lebesgue spaces $L^{p}$ for which the norm of the Riesz projection $P$ is equal to one. On the other hand, [5, Theorem 7.7] implies that

$$
\|S\|_{H^{2} \rightarrow H^{2}}=1, \quad\|S\|_{H^{p} \rightarrow H^{p}}>1 \quad \text { for } \quad p \in(1, \infty) \backslash\{2\} .
$$

Thus, for the class of Lebesgue spaces $L^{p}$ with $1<p<\infty$, one has

$$
\begin{equation*}
\|P\|_{L^{p} \rightarrow L^{p}}=1 \quad \Longleftrightarrow \quad\|S\|_{H^{p} \rightarrow H^{p}}=1 \tag{1.2}
\end{equation*}
$$

Banach function spaces provide a far-reaching generalisation of Lebesgue spaces. The class of Banach function spaces includes Lebesgue spaces $L^{p}, 1 \leq p \leq \infty$, Orlicz spaces $L^{\varphi}$, Lorentz spaces $L^{p, q}$, all other rearrangement-invariant spaces (see, e.g., [3, Ch. 2 and 4]), as well as, variable Lebesgue spaces $L^{p(\cdot)}$ (see, e.g., [10]), which are not rearrangement-invariant.

It follows from [21, Theorem 4.5, Corollary 4.6] that if $X$ is a reflexive rearrangementinvariant Banach function space such that $P: X \rightarrow X$ is bounded, then

$$
\begin{align*}
\|P\|_{X \rightarrow X}^{\text {ess }} & :=\inf \left\{\|P-K\|_{X \rightarrow X}: K \text { is compact on } X\right\} \\
& \geq \frac{1}{\sin \left(\pi \min \left\{p_{X}, 1-q_{X}\right\}\right)} \tag{1.3}
\end{align*}
$$

where $p_{X}$ and $q_{X}$ are the Zippin indices of the space $X$ (see [28, pp. 27-28] for their definition and the proof of the inequalities $0 \leq p_{X} \leq q_{X} \leq 1$, which are valid for arbitrary rearrangement-invariant Banach function spaces).

So, if $p_{X} \neq 1 / 2$ or $q_{X} \neq 1 / 2$, then $\|P\|_{X \rightarrow X} \geq\|P\|_{X \rightarrow X}^{\text {ess }}>1$. We note in passing that if $X$ is a rearrangement-invariant Banach function space such that $P: X \rightarrow X$ is bounded, then $P$ is maximally noncompact on $X$, that is,

$$
\|P\|_{X \rightarrow X}=\|P\|_{X \rightarrow X}^{\operatorname{ess}^{s .}}
$$

(see [24, Theorem 1.1]).
Estimate (1.3) does not exclude the possibility of $\|P\|_{X \rightarrow X}=1$ if $p_{X}=q_{X}=1 / 2$. Note that, for instance, the Lorentz spaces $L^{2, r}, 1<r<\infty$, are reflexive rearrangementinvariant Banach function spaces (see, e.g., [3, Ch. 4, Section 4]) with the Zippin indices $p_{L^{2, r}}=q_{L^{2, r}}=1 / 2$ (see, e.g., [28, pp. 27-28]), and the operator $P$ is bounded on $L^{2, r}$ for every $1<r<\infty$ (the latter follows from Calderón's extension of the Marcinkiewicz interpolation theorem [3, Ch. 4, Theorem 4.13]). On the other hand, it follows from Holmstedt's formula (see [19, Theorems 4.2-4.3]) for the $K$-functional for Lorentz spaces that for $\delta \in(0,1)$ and $1 \leq r \leq \infty$, the space $X_{\delta, r}:=\left(L^{2 /(1-\delta)}, L^{2 /(1+\delta)}\right)_{1 / 2, r}$, obtained from the Lebesgue spaces $L^{2 /(1-\delta)}$ and $L^{2 /(1+\delta)}$ by the $K$-method of real interpolation (see, e.g., [3, Ch. 5]), coincides with the Lorentz space $L^{2, r}$ up to equivalence of the norms. Since the $K$-method of real interpolation is exact (see, e.g., [3, Ch. 5, Theorem 1.12]), we conclude from (1.1) that

$$
\begin{aligned}
\|P\|_{X_{\delta, r} \rightarrow X_{\delta, r}} & \leq\|P\|_{L^{2 /(1-\delta)} \rightarrow L^{2 /(1-\delta)}}^{1 / 2}\|P\|_{L^{2 /(1+\delta)} \rightarrow L^{2 /(1+\delta)}}^{1 / 2} \\
& =\frac{1}{\sqrt{\sin \frac{\pi(1-\delta)}{2}}} \frac{1}{\sqrt{\sin \frac{\pi(1+\delta)}{2}}}=\frac{1}{\sin \frac{\pi(1+\delta)}{2}} .
\end{aligned}
$$

Hence, for every $\varepsilon>0$ and $r \in[1, \infty]$, one can find $\delta>0$ such that

$$
\begin{equation*}
\|P\|_{X_{\delta, r} \rightarrow X_{\delta, r}} \leq \frac{1}{\sin \frac{\pi(1+\delta)}{2}}<1+\varepsilon \tag{1.4}
\end{equation*}
$$

Thus, for every $\varepsilon>0$ and $r \in[1, \infty]$, one can equip the Lorentz space $L^{2, r}$ with an equivalent norm $\|\cdot\|_{L_{\varepsilon}^{2, r}}$ such that $\|P\|_{L_{\varepsilon}^{2, r} \rightarrow L_{\varepsilon}^{2, r}}<1+\varepsilon$ (it is enough to take $\|\cdot\|_{L_{\varepsilon}^{2, r}}:=$ $\|\cdot\|_{X_{\delta, r}}$, where $\delta$ satisfies (1.4)).

So, the following natural question arises: can the norm of the Riesz projection $P$ on a (not necessarily rearrangement-invariant) Banach function space $X$ be equal to one if $X$ does not coincide with $L^{2}$ ? The first main result of the paper gives a negative answer to this question.

Theorem 1.1 (Main result 1). Let $X$ be a Banach function space such that $\|P\|_{X \rightarrow X}=1$. Then $X$ coincides with $L^{2}$ and there exists a constant $C \in(0, \infty)$ such that

$$
\begin{equation*}
\|g\|_{X}=C\|g\|_{L^{2}} \quad \text { for all } \quad g \in X \tag{1.5}
\end{equation*}
$$

Our second main result deals with the extension of property (1.2) to the setting of Banach function spaces.

Theorem 1.2 (Main result 2). Let $X$ be a Banach function space. Then $\|P\|_{X \rightarrow X}=1$ if and only if $\|S\|_{H[X] \rightarrow H[X]}=1$.

The paper is organized as follows. In Section 2, we recall definitions of a Banach function space and its associate space $X^{\prime}$, of the subspace $X_{a}$ of all functions of absolutely continuous norm and of the subspace $X_{b}$, which is the closure of the set of all simple functions in $X$. Further, we note that if $X_{a}=X_{b}$, then the set of trigonometric polynomials $\mathcal{P}$ is dense in $X_{b}$. We also need a few notions from the theory of analytic functions on the open unit disk $\mathbb{D}$, Poisson integrals, and the Hilbert transform $H$. After these preliminaries, we recall that if $f \in L^{p}$ is a real-valued function and $u$ is an inner function vanishing at zero, then $H(f \circ u)=(H f) \circ u$. We conclude Section 2 by recalling several known facts about the Riesz projection scattered in our previous papers. We start Section 3 by proving a property of the norm in a real Hilbert space, and then give a proof of Theorem 1.1.

As far as Theorem 1.2 is concerned, the proof of the implications

$$
\begin{align*}
& \|P\|_{X \rightarrow X}=1 \quad \Longrightarrow \quad\|S\|_{H[X] \rightarrow H[X]}=1 \\
& \|S\|_{H[X] \rightarrow H[X]}=1 \quad \Longrightarrow \quad\|P g\|_{X} \leq\|g\|_{X} \quad \text { for all continuous } g \tag{1.6}
\end{align*}
$$

is not difficult. The main difficulty lies in extending the estimate $\|P g\|_{X} \leq\|g\|_{X}$ in (1.6) to all $g \in X$ when $X$ is not separable. This difficulty is addressed in Section 4 where we refine [23, Theorem 3.7] and [25, Theorem 3.3] and show that if the Hilbert transform
$H$ is of weak type from the space $C$ of continuous functions to a Banach function space $X$, then $X_{a}=X_{b}$. This implies that if the Riesz projection $P$ is bounded from $C$ to a Banach function space $X$, then $X_{a}=X_{b}$ and $\mathcal{P}$ is dense in $X_{b}$. This observation is a key ingredient in the proof of Theorem 1.2 given in Section 5.

Finally, in Section 6, we extend [5, Theorem 7.7] to the setting of Banach function spaces $X$ and show that the norm of $P$ on $X$ can be expressed in terms of Toeplitz operators.

## 2. Preliminaries

### 2.1. Banach function spaces and their associate spaces

Let $\mathcal{M}$ be the set of all measurable extended complex-valued functions on $\mathbb{T}$ equipped with the normalized measure $d m(t)=|d t| /(2 \pi)$ and let $\mathcal{M}^{+}$be the subset of functions in $\mathcal{M}$ whose values lie in $[0, \infty]$.

A mapping $\rho: \mathcal{M}^{+} \rightarrow[0, \infty]$ is called a Banach function norm if, for all functions $f, g, f_{n} \in \mathcal{M}^{+}$with $n \in \mathbb{N}$, and for all constants $a \geq 0$, the following properties hold:
(A4) $\rho(1)<\infty$,

$$
\begin{align*}
& \text { (A1) } \\
& \rho(f)=0 \Leftrightarrow f=0 \text { a.e., } \rho(a f)=a \rho(f), \rho(f+g) \leq \rho(f)+\rho(g), \\
& \text { (A2) }  \tag{A1}\\
& 0 \leq g \leq f \text { a.e. } \Rightarrow \rho(g) \leq \rho(f) \quad \text { (the lattice property), } \\
& \text { (A3) } \\
& 0 \leq f_{n} \uparrow f \text { a.e. } \Rightarrow \rho\left(f_{n}\right) \uparrow \rho(f) \quad \text { (the Fatou property), }  \tag{A5}\\
& \text { (A4) } \\
& \rho(1)<\infty, \\
& \text { (A5) } \\
& \int_{\mathbb{T}} f(t) d m(t) \leq C \rho(f)
\end{align*}
$$

with a constant $C \in(0, \infty)$ that may depend on $\rho$, but is independent of $f$. When functions differing only on a set of measure zero are identified, the set $X$ of all functions $f \in \mathcal{M}$ for which $\rho(|f|)<\infty$ is called a Banach function space. For each $f \in X$, the norm of $f$ is defined by $\|f\|_{X}:=\rho(|f|)$. The set $X$ equipped with the natural linear space operations and this norm becomes a Banach space. If $\rho$ is a Banach function norm, its associate norm $\rho^{\prime}$ defined on $\mathcal{M}^{+}$by

$$
\rho^{\prime}(g):=\sup \left\{\int_{\mathbb{T}} f(t) g(t) d m(t): f \in \mathcal{M}^{+}, \rho(f) \leq 1\right\}, \quad g \in \mathcal{M}^{+}
$$

is a Banach function norm itself. The Banach function space $X^{\prime}$ determined by the Banach function norm $\rho^{\prime}$ is called the associate space (Köthe dual) of $X$. The associate space $X^{\prime}$ can be viewed as a subspace of the Banach dual space $X^{*}$ (see [3, Ch. 1, Sections 1-2]).

### 2.2. Density of trigonometric polynomials in the subspace $X_{b}$

The characteristic (indicator) function of a measurable set $E \subset \mathbb{T}$ is denoted by $\chi_{E}$. A function $f$ in a Banach function space $X$ is said to have absolutely continuous norm in $X$ if $\left\|f \chi_{\gamma_{n}}\right\|_{X} \rightarrow 0$ for every sequence $\left\{\gamma_{n}\right\}$ of measurable sets such that $\chi_{\gamma_{n}} \rightarrow 0$ almost everywhere as $n \rightarrow \infty$. The set of all functions of absolutely continuous norm in $X$ is denoted by $X_{a}$. If $X_{a}=X$, then one says that $X$ has absolutely continuous norm. Let $S_{0}$ be the set of all simple functions on $\mathbb{T}$. Let $X_{b}$ denote the closure of $S_{0}$ in the norm of $X$. We refer to [3, Ch. 1, Section 3] for properties of the subspaces $X_{a}$ and $X_{b}$.

For $n \in \mathbb{Z}_{+}:=\{0,1,2, \ldots\}$, a function of the form $\sum_{k=-n}^{n} \alpha_{k} \mathbf{e}_{k}$, where $\alpha_{k} \in \mathbb{C}$ for all $k \in\{-n, \ldots, n\}$, is called a trigonometric polynomial of order $n$. The set of all trigonometric polynomials is denoted by $\mathcal{P}$.

Lemma 2.1 ([23, Lemma 2.1]). Let $X$ be a Banach function space. If $X_{a}=X_{b}$, then the set of trigonometric polynomials $\mathcal{P}$ is dense in $X_{b}$.

### 2.3. Classes of analytic functions on the open unit disk

Let $\mathbb{D}$ denote the open unit disk in the complex plane $\mathbb{C}$. Recall that a function $F$ analytic in $\mathbb{D}$ is said to belong to the Hardy space $H^{p}(\mathbb{D}), 0<p \leq \infty$, if

$$
\begin{aligned}
\|F\|_{H^{p}(\mathbb{D})} & :=\sup _{0 \leq r<1}\left(\frac{1}{2 \pi} \int_{-\pi}^{\pi}\left|F\left(r e^{i \theta}\right)\right|^{p} d \theta\right)^{1 / p}<\infty, \quad 0<p<\infty \\
\|F\|_{H^{\infty}(\mathbb{D})} & :=\sup _{z \in \mathbb{D}}|F(z)|<\infty
\end{aligned}
$$

Let $g$ be a measurable function on $\mathbb{T}$ with $\log |g| \in L^{1}$. An outer function (of absolute value $|g|$ ) is a function $f=\lambda G$ with $\lambda \in \mathbb{C},|\lambda|=1$, and

$$
G(z):=\exp \left(\frac{1}{2 \pi} \int_{-\pi}^{\pi} \frac{e^{i \theta}+z}{e^{, \theta}-z} \log \left|g\left(e^{i \theta}\right)\right| d \theta\right), \quad z \in \mathbb{D} .
$$

The Smirnov class $\mathcal{D}(\mathbb{D})$ consists of all functions $f$ analytic in $\mathbb{D}$, which can be represented in the form $f=f_{1} / f_{2}$, where $f_{2}$ is outer and $f_{1}, f_{2} \in \bigcup_{0<p \leq \infty} H^{p}(\mathbb{D})$ (see, e.g., [29, Definition 3.3.1]). Recall that an inner function is a function $u \in H^{\infty}(\mathbb{D})$ such that $\left|u\left(e^{i \theta}\right)\right|=1$ for a.e. $\theta \in[-\pi, \pi]$.

Lemma 2.2. If $u$ is an inner function such that $u(0)=0$, then $u$ is a measure preserving transformation from $\mathbb{T}$ onto itself.

This lemma goes back to Nordgren (see corollary to [30, Lemma 1] and also [9, Remark 9.4.6], [23, Lemma 2.5], [12, Theorem 5.5]).

For a finite collection $z_{1}, \ldots, z_{n} \in \mathbb{D}$ and $\gamma \in \mathbb{T}$, the function

$$
B(z)=\gamma \prod_{j=1}^{n} \frac{z-z_{j}}{1-\overline{z_{j}} z}
$$

is called a finite Blaschke product. As is well known, every finite Blaschke product satisfies

$$
|B(z)|<1 \quad \text { for } \quad z \in \mathbb{D}, \quad|B(\zeta)|=1 \quad \text { for } \quad \zeta \in \mathbb{T}
$$

(see, e.g., [13, Section 3.1]).

### 2.4. The Hilbert transform and Poisson integrals

The Hilbert transform of a function $f \in L^{1}$ is defined by

$$
(H f)\left(e^{i \vartheta}\right):=\frac{1}{2 \pi} \text { p.v. } \int_{-\pi}^{\pi} f\left(e^{i \theta}\right) \cot \frac{\vartheta-\theta}{2} d \theta, \quad \vartheta \in[-\pi, \pi] .
$$

For $\vartheta \in[-\pi, \pi]$ and $r \in[0,1)$, let

$$
P_{r}(\vartheta):=\frac{1-r^{2}}{1-2 r \cos \vartheta+r^{2}}, \quad Q_{r}(\vartheta):=\frac{2 r \sin \vartheta}{1-2 r \cos \vartheta+r^{2}}
$$

be the Poisson kernel and the conjugate Poisson kernel, respectively.
Theorem 2.3. Let $1<p<\infty$. If $f \in L^{p}$ is a real-valued function, then the function defined by

$$
\begin{equation*}
w\left(r e^{i \vartheta}\right)=\frac{1}{2 \pi} \int_{-\pi}^{\pi} f\left(e^{i \theta}\right)\left(P_{r}+i Q_{r}\right)(\vartheta-\theta) d \theta, \vartheta \in[-\pi, \pi], r \in[0,1) \tag{2.1}
\end{equation*}
$$

belongs to the Hardy space $H^{p}(\mathbb{D})$ and $\operatorname{Im} w(0)=0$. Its nontangential boundary values $w\left(e^{i \vartheta}\right)$ as $z \rightarrow e^{i \vartheta}$ exist for a.e. $\vartheta \in[-\pi, \pi]$ and

$$
\begin{equation*}
\operatorname{Re} w\left(e^{i \vartheta}\right)=f\left(e^{i \vartheta}\right), \operatorname{Im} w\left(e^{i \vartheta}\right)=(H f)\left(e^{i \vartheta}\right) \text { for a.e. } \vartheta \in[-\pi, \pi] . \tag{2.2}
\end{equation*}
$$

This statement is well known (see, e.g., [27, Ch. I, Section D and Ch. V, Section B. $\left.2^{\circ}\right]$ ).

The next lemma will play an important role in the proof of Theorem 1.2.
Lemma 2.4. Let $1<p<\infty, f \in L^{p}$ be a real-valued function and $u$ be an inner function such that $u(0)=0$. Then $H(f \circ u)=(H f) \circ u$.

Proof. By Lemma 2.2, $u: \mathbb{T} \rightarrow \mathbb{T}$ is a measure preserving transformation. Therefore the operator $g \mapsto g \circ u$ is isometric, and hence bounded, on $L^{p}$. So, it is sufficient to prove the equality $H(f \circ u)=(H f) \circ u$ for all $f$ from a dense subset of $L^{p}$.

We will suppose that $f$ is Hölder continuous. Then it follows from [17, Ch. IX, §1, Theorem 1] that

$$
F\left(r e^{i \vartheta}\right):=\frac{1}{2 \pi} \int_{-\pi}^{\pi} f\left(e^{i \theta}\right) P_{r}(\vartheta-\theta) d \theta, \quad \vartheta \in[-\pi, \pi], \quad r \in[0,1]
$$

is continuous in $\overline{\mathbb{D}}$ and

$$
\begin{equation*}
\operatorname{Re} w\left(e^{i \vartheta}\right)=F\left(e^{i \vartheta}\right)=f\left(e^{i \vartheta}\right) \quad \text { for all } \quad \vartheta \in[-\pi, \pi], \tag{2.3}
\end{equation*}
$$

where $w$ is the function defined by (2.1). Further, [17, Ch. IX, §5, Theorem 5] implies that $w$ is Hölder continuous on $\overline{\mathbb{D}}$. Hence $\operatorname{Im} w\left(e^{i \vartheta}\right)$ is Hölder continuous on $[-\pi, \pi]$. It follows from Theorem 2.3 that $\operatorname{Im} w\left(e^{i \vartheta}\right)=(H f)\left(e^{i \vartheta}\right)$ for a.e. $\vartheta \in[-\pi, \pi]$. Since $f$ is Hölder continuous, we conclude from Privalov's theorem (see, e.g., [9, Theorem 3.1.1] or [33, Ch. III, Theorem 13.29]) that $H f$ is Hölder continuous. Since the functions $\operatorname{Im} w\left(e^{i \vartheta}\right)$ and $(H f)\left(e^{i \vartheta}\right)$ are equal almost everywhere and they are continuous, we conclude that they are equal everywhere:

$$
\begin{equation*}
\operatorname{Im} w\left(e^{i \vartheta}\right)=(H f)\left(e^{i \vartheta}\right) \quad \text { for all } \quad \vartheta \in[-\pi, \pi] . \tag{2.4}
\end{equation*}
$$

Let $W:=w \circ u$. Then $W \in H^{p}(\mathbb{D})($ see [11, Section 2.6, Corollary to Theorem 2.12])), and $\operatorname{Im} W(0)=\operatorname{Im} w(u(0))=\operatorname{Im} w(0)=0$. It follows from (2.3)-(2.4) and $u\left(e^{i \vartheta}\right) \in \mathbb{T}$ for a.e. $\vartheta \in[-\pi, \pi]$ that

$$
\begin{align*}
& \operatorname{Re} W\left(e^{i \vartheta}\right)=\operatorname{Re}(w \circ u)\left(e^{i \vartheta}\right)=(f \circ u)\left(e^{i \vartheta}\right),  \tag{2.5}\\
& \operatorname{Im} W\left(e^{i \vartheta}\right)=\operatorname{Im}(w \circ u)\left(e^{i \vartheta}\right)=((\operatorname{Im} w) \circ u)\left(e^{i \vartheta}\right)=((H f) \circ u)\left(e^{i \vartheta}\right) \tag{2.6}
\end{align*}
$$

for a.e. $\vartheta \in[-\pi, \pi]$. According to Theorem 2.3,

$$
\begin{equation*}
\operatorname{Im} W\left(e^{i \vartheta}\right)=(H(\operatorname{Re} W))\left(e^{i \vartheta}\right) \tag{2.7}
\end{equation*}
$$

for a.e. $\vartheta \in[-\pi, \pi]$ (see (2.2)). Combining (2.5)-(2.7), we get

$$
((H f) \circ u)\left(e^{i \vartheta}\right)=\operatorname{Im} W\left(e^{i \vartheta}\right)=(H(\operatorname{Re} W))\left(e^{i \vartheta}\right)=(H(f \circ u))\left(e^{i \vartheta}\right)
$$

for a.e. $\vartheta \in[-\pi, \pi]$.

### 2.5. Some known facts on the Riesz projection

In this subsection we list several known facts about the operator $P$, which will be used in this paper.

Lemma 2.5 ([23, formula (1.4)]). If $f \in L^{1}$ is such that $P f \in L^{1}$, then

$$
(P f)^{\wedge}(n)= \begin{cases}\widehat{f}(n), & \text { if } \quad n \geq 0 \\ 0, & \text { if } \quad n<0\end{cases}
$$

Lemma 2.6 ([22, Lemma 3.1]). Let $f \in L^{1}$. Suppose there exists $g \in H^{1}$ such that $\widehat{f}(n)=\widehat{g}(n)$ for all $n \geq 0$. Then $P f=g$.

Theorem 2.7 ([25, Theorem 3.4]). Let $X$ be a Banach function space and $X^{\prime}$ be its associate space. If $P: X_{b} \rightarrow X$ is bounded, then $P: X \rightarrow X$ is bounded, $P$ maps $X_{b}$ into itself,

$$
\begin{equation*}
\|P\|_{X \rightarrow X}=\|P\|_{X_{b} \rightarrow X_{b}}, \tag{2.8}
\end{equation*}
$$

and the adjoint of the bounded operator $P: X_{b} \rightarrow X_{b}$ is the operator $P: X^{\prime} \rightarrow X^{\prime}$, which implies that the latter is also bounded.

## 3. Proof of the first main result

### 3.1. A property of the norm of a real Hilbert space

Lemma 3.1. Let $\mathcal{H}$ be a real Hilbert space, $\varrho$ be a norm equivalent to $\|\cdot\|_{\mathcal{H}}$, and $\varrho^{\prime}$ be the associate norm,

$$
\varrho^{\prime}(x):=\sup \left\{\left|\langle y, x\rangle_{\mathcal{H}}\right|: y \in \mathcal{H}, \varrho(y) \leq 1\right\}, \quad x \in \mathcal{H},
$$

where $\langle\cdot, \cdot\rangle_{\mathcal{H}}$ denotes the inner product in $\mathcal{H}$. If

$$
\begin{equation*}
\varrho(x) \varrho^{\prime}(x)=\|x\|_{\mathcal{H}}^{2} \quad \text { for all } \quad x \in \mathcal{H} \tag{3.1}
\end{equation*}
$$

then there exists a constant $C \in(0, \infty)$ such that

$$
\begin{equation*}
\varrho(x)=C\|x\|_{\mathcal{H}} \quad \text { for all } \quad x \in \mathcal{H} \tag{3.2}
\end{equation*}
$$

Proof. Fix $a \in \mathcal{H} \backslash\{0\}$ and put

$$
\begin{equation*}
C:=\frac{\varrho(a)}{\|a\|_{\mathcal{H}}} . \tag{3.3}
\end{equation*}
$$

If $x=0$, then (3.2) holds trivially. If $x \in \mathcal{H} \backslash\{0\}$ and $a$ are linearly dependent, then there exists $\lambda \in \mathbb{C} \backslash\{0\}$ such that $x=\lambda a$, and

$$
C=\frac{|\lambda| \varrho(a)}{|\lambda|\|a\|_{\mathcal{H}}}=\frac{\varrho(\lambda a)}{\|\lambda a\|_{\mathcal{H}}}=\frac{\varrho(x)}{\|x\|_{\mathcal{H}}},
$$

which implies (3.2).
Now suppose that $a$ and $x \in \mathcal{H} \backslash\{0\}$ are linearly independent. Let $\mathcal{L}$ be the twodimensional subspace spanned by $a$ and $x$. Choosing an orthonormal basis in $\mathcal{L}$ we can identify $\mathcal{L}$ with $\mathbb{R}^{2}$ and $\|\cdot\|_{\mathcal{H}}$ with the standard Euclidean norm $\|\cdot\|$ on $\mathbb{R}^{2}$. With a slight abuse of notation, we denote the norms generated by $\varrho$ and $\varrho^{\prime}$ on $\mathbb{R}^{2}$ by the same symbols.

Since $\varrho$ is positively homogeneous of degree 1, it can be represented in the form

$$
\begin{equation*}
\varrho(r \cos \theta, r \sin \theta)=r \Phi(\theta), \quad r>0, \quad \theta \in[0,2 \pi) \tag{3.4}
\end{equation*}
$$

where

$$
\Phi(\theta):=\varrho(\cos \theta, \sin \theta), \quad \theta \in[0,2 \pi) .
$$

Let

$$
m:=\inf _{\theta \in[0,2 \pi)} \Phi(\theta) .
$$

Then

$$
\begin{equation*}
\Phi(\theta) \geq m>0 \quad \text { for all } \quad \theta \in[0,2 \pi) . \tag{3.5}
\end{equation*}
$$

On the other hand, since all norms on $\mathbb{R}^{2}$ are equivalent, there exists $M \in(0, \infty)$ such that $\varrho(\cdot) \leq M\|\cdot\|$. Then

$$
\begin{align*}
\left|\Phi(\theta)-\Phi\left(\theta^{\prime}\right)\right| & =\left|\varrho(\cos \theta, \sin \theta)-\varrho\left(\cos \theta^{\prime}, \sin \theta^{\prime}\right)\right| \\
& \leq \varrho\left((\cos \theta, \sin \theta)-\left(\cos \theta^{\prime}, \sin \theta^{\prime}\right)\right) \\
& \leq M\left\|(\cos \theta, \sin \theta)-\left(\cos \theta^{\prime}, \sin \theta^{\prime}\right)\right\| \\
& =2 M\left|\sin \frac{\theta-\theta^{\prime}}{2}\right| \leq M\left|\theta-\theta^{\prime}\right| \tag{3.6}
\end{align*}
$$

for all $\theta, \theta^{\prime} \in[0,2 \pi)$. It follows from (3.5)-(3.6) that $R:=1 / \Phi$ is also Lipschitz continuous and hence absolutely continuous.

Take any

$$
w \in S_{\varrho}:=\left\{z \in \mathbb{R}^{2}: \varrho(z)=1\right\}=\{(R(\theta) \cos \theta, R(\theta) \sin \theta): \theta \in[0,2 \pi)\}
$$

It follows from (3.1) and the definition of $\varrho^{\prime}$ that the function

$$
S_{\varrho} \ni z \mapsto\langle z, w\rangle
$$

achieves its maximum on $S_{\varrho}$ at $z=w$. In other words, for any $\theta_{0} \in[0,2 \pi)$, the function

$$
\begin{aligned}
F(\theta) & :=\left\langle(R(\theta) \cos \theta, R(\theta) \sin \theta),\left(R\left(\theta_{0}\right) \cos \theta_{0}, R\left(\theta_{0}\right) \sin \theta_{0}\right)\right\rangle_{\mathcal{H}} \\
& =R\left(\theta_{0}\right) R(\theta)\left(\cos \theta \cos \theta_{0}+\sin \theta \sin \theta_{0}\right)=R\left(\theta_{0}\right) R(\theta) \cos \left(\theta-\theta_{0}\right)
\end{aligned}
$$

achieves its maximum at $\theta=\theta_{0}$. If $R$ is differentiable at $\theta_{0}$, then

$$
0=F^{\prime}\left(\theta_{0}\right)=\left.R\left(\theta_{0}\right)\left(R^{\prime}(\theta) \cos \left(\theta-\theta_{0}\right)-R(\theta) \sin \left(\theta-\theta_{0}\right)\right)\right|_{\theta=\theta_{0}}=R\left(\theta_{0}\right) R^{\prime}\left(\theta_{0}\right)
$$

Hence $R$ is an absolutely continuous function with $R^{\prime}=0$ a.e. So, $R$ is constant. Then it follows from (3.4) that $\varrho(z) /\|z\|$ is constant for $z \in \mathcal{L} \backslash\{0\}$. This observation and (3.3) imply that

$$
C=\frac{\varrho(a)}{\|a\|_{\mathcal{H}}}=\frac{\varrho(a)}{\|a\|}=\frac{\varrho(x)}{\|x\|}=\frac{\varrho(x)}{\|x\|_{\mathcal{H}}}
$$

which implies (3.2) in the case when $x$ and $a$ are linearly independent.

### 3.2. Proof of Theorem 1.1

Since $P: X \rightarrow X$ is bounded, we have $X_{a}=X_{b}$ (see [23, Theorem 3.7]) and $\left(X_{b}\right)^{*}=$ $X^{\prime}$ (see [3, Ch. 1, Corollary 4.2]). Take any $\varepsilon>0$ and any $g \in X_{b}$ such that $|g| \geq \varepsilon$ a.e. on $\mathbb{T}$. Put $\log ^{+}|z|:=\max \{0, \log |z|\}$ for $z \in \mathbb{C}$. Since

$$
\log |z|=\log ^{+}|z|-\log ^{+}(1 /|z|), \quad \log ^{+}|z| \leq|z|, \quad z \in \mathbb{C}
$$

it follows from $|g| \geq \varepsilon$ a.e. on $\mathbb{T}$ and $g \in L^{1}$ that $\log |g| \in L^{1}$. Then, by Szegö's theorem (see, e.g., [29, Theorem 2.6.1]), the outer function

$$
G(z):=\exp \left(\frac{1}{2 \pi} \int_{-\pi}^{\pi} \frac{e^{i \theta}+z}{e^{-\theta}-z} \log \left|g\left(e^{i \theta}\right)\right| d \theta\right), \quad z \in \mathbb{D},
$$

belongs to $H^{1}(\mathbb{D})$ and $|G|=|g|$ a.e. on $\mathbb{T}$. Then $G \in H^{1} \cap X=H[X]$. Since $P$ is bounded on $X$, it follows from Lemma 2.5 and the uniqueness theorem for Fourier series (see, e.g., [26, Ch. 1, Theorem 2.7]) that $P G=G$ a.e. on $\mathbb{T}$. Taking into account that $g \in X_{b}$ and $|G|=|g|$, we deduce from [3, Ch. 1, Theorem 3.11 and Definition 3.7]) that, in fact, $G \in X_{b}$. By the Hahn-Banach theorem, there exists $\varphi \in\left(X_{b}\right)^{*}$ such that $\|\varphi\|_{\left(X_{b}\right)^{*}}=1$ and $\varphi(G)=\|G\|_{X_{b}}$. Since $\left(X_{b}\right)^{*}$ is isometrically isomorphic to $X^{\prime}$, there exists $u \in X^{\prime}$ such that $\|\varphi\|_{\left(X_{b}\right)^{*}}=\|u\|_{X^{\prime}}$ and

$$
\varphi(f)=\frac{1}{2 \pi} \int_{-\pi}^{\pi} f\left(e^{i \theta}\right) \overline{u\left(e^{i \theta}\right)} d \theta, \quad f \in X_{b}
$$

Thus $\|u\|_{X^{\prime}}=1$ and

$$
\|G\|_{X}=\|G\|_{X_{b}}=\frac{1}{2 \pi} \int_{-\pi}^{\pi} G\left(e^{i \theta}\right) \overline{u\left(e^{i \theta}\right)} d \theta
$$

It follows from Theorem 2.7 that the adjoint operator of $P: X_{b} \rightarrow X_{b}$ is the operator $P^{*}=P: X^{\prime} \rightarrow X^{\prime}$ and

$$
\|P\|_{X^{\prime} \rightarrow X^{\prime}}=\|P\|_{X_{b} \rightarrow X_{b}}=\|P\|_{X \rightarrow X}=1
$$

So, the function

$$
u_{+}:=P u \in H\left[X^{\prime}\right] \subset H^{1}
$$

satisfies $\left\|u_{+}\right\|_{X^{\prime}} \leq 1$ and

$$
\begin{align*}
\|G\|_{X} & =\frac{1}{2 \pi} \int_{-\pi}^{\pi} G\left(e^{i \theta}\right) \overline{u\left(e^{i \theta}\right)} d \theta=\frac{1}{2 \pi} \int_{-\pi}^{\pi}(P G)\left(e^{i \theta}\right) \overline{u\left(e^{i \theta}\right)} d \theta \\
& =\frac{1}{2 \pi} \int_{-\pi}^{\pi} G\left(e^{i \theta}\right) \overline{(P u)\left(e^{i \theta}\right)} d \theta=\frac{1}{2 \pi} \int_{-\pi}^{\pi} G\left(e^{i \theta}\right) \overline{u_{+}\left(e^{i \theta}\right)} d \theta \tag{3.7}
\end{align*}
$$

(see [25, Lemma 4.1]). Using Hölder's inequality (see [3, Ch. 1, Theorem 2.4]) and taking into account that $\left\|u_{+}\right\|_{X^{\prime}} \leq 1$, one gets

$$
\begin{aligned}
\|G\|_{X} & =\operatorname{Re}\left(\frac{1}{2 \pi} \int_{-\pi}^{\pi} G\left(e^{i \theta}\right) \overline{u_{+}\left(e^{i \theta}\right)} d \theta\right)=\frac{1}{2 \pi} \int_{-\pi}^{\pi} \operatorname{Re}\left(G\left(e^{i \theta}\right) \overline{u_{+}\left(e^{i \theta}\right)}\right) d \theta \\
& \leq \frac{1}{2 \pi} \int_{-\pi}^{\pi}\left|G\left(e^{i \theta}\right)\right|\left|u_{+}\left(e^{i \theta}\right)\right| d \theta \leq\|G\|_{X}\left\|u_{+}\right\|_{X^{\prime}} \leq\|G\|_{X}
\end{aligned}
$$

Then $\left\|u_{+}\right\|_{X^{\prime}}=1$,

$$
\begin{equation*}
\frac{1}{2 \pi} \int_{-\pi}^{\pi} \operatorname{Re}\left(G\left(e^{i \theta}\right) \overline{u_{+}\left(e^{i \theta}\right)}\right) d \theta=\frac{1}{2 \pi} \int_{-\pi}^{\pi}\left|G\left(e^{i \theta}\right)\right|\left|u_{+}\left(e^{i \theta}\right)\right| d \theta \tag{3.8}
\end{equation*}
$$

Since $|G|\left|u_{+}\right|-\operatorname{Re}\left(G \overline{u_{+}}\right) \geq 0$ a.e. on $\mathbb{T}$, it follows from (3.8) that

$$
\operatorname{Re}\left(G \overline{u_{+}}\right)=|G|\left|u_{+}\right| \quad \text { a.e. on } \quad \mathbb{T} .
$$

Hence $\psi:=G \overline{u_{+}} \in L^{1}$ is a nonnegative function and $\log \psi=\log |G|+\log \left|u_{+}\right|$. Since $|G|=|g| \geq \varepsilon$ a.e. on $\mathbb{T}$ and $\left\|u_{+}\right\|_{X^{\prime}}=1$, we have $G, u_{+} \neq 0$. Taking into account that $G, u_{+} \in H^{1}$, we deduce from [29, Corollary 2.2.3] that $\log |G| \in L^{1}$ and $\log \left|u_{+}\right| \in L^{1}$. Thus $\log \psi \in L^{1}$. By Szegő's theorem (see, e.g., [29, Theorem 2.6.1]), there exists an outer function $\Psi \in H^{2}(\mathbb{D})$ such that $|\Psi|=\psi^{1 / 2}$ a.e. on $\mathbb{T}$. So,

$$
G \overline{u_{+}}=\psi=|\Psi|^{2}=\Psi \bar{\Psi} \quad \text { a.e. on } \quad \mathbb{T},
$$

whence

$$
\begin{equation*}
\overline{\left(\frac{u_{+}}{\Psi}\right)}=\frac{\Psi}{G} \quad \text { a.e. on } \quad \mathbb{T} . \tag{3.9}
\end{equation*}
$$

Since $\Psi \in H^{2}(\mathbb{D})$ and $G \in H^{1}(\mathbb{D})$ are outer functions, we conclude that $\Psi / G$ belongs to the Smirnov class $\mathcal{D}(\mathbb{D})$. Moreover, $\Psi \in L^{2}$ and $1 / G \in L^{\infty}$. Hence $\Psi / G \in L^{2}$. Then, in view of a generalization of Smirnov's theorem (see, e.g., [29, Section 3.3.1(g)] or [11, Theorem 2.11]), $\Psi / G \in H^{2}(\mathbb{D}) \subset H^{1}(\mathbb{D})$. Similarly, $\Psi \in H^{2}(\mathbb{D})$ is an outer function and $u_{+} \in H\left[X^{\prime}\right] \subset H^{1}$. Let us extend $u_{+}$to the unit disk analytically:

$$
u_{+}(z):=\frac{1}{2 \pi} \int_{-\pi}^{\pi} u_{+}\left(e^{i \theta}\right) \frac{1-r^{2}}{1-2 r \cos (\varphi-\theta)+r^{2}} d \theta, \quad z=r e^{i \varphi} \in \mathbb{D} .
$$

Then $u_{+} \in H^{1}(\mathbb{D})$. So, $u_{+} / \Psi \in \mathcal{D}(\mathbb{D})$. On the other hand,

$$
\frac{u_{+}}{\Psi}=\overline{\left(\frac{\Psi}{G}\right)} \in L^{2} .
$$

Hence, applying Smirnov's theorem once again, one gets

$$
\frac{u_{+}}{\Psi} \in H^{2}(\mathbb{D}) \subset H^{1}(\mathbb{D})
$$

So, we have shown that $F:=u_{+} / \Psi \in H^{2}(\mathbb{D})$ and $\bar{F} \in H^{2}(\mathbb{D})$ (see (3.9)). Taking into account that $\overline{\widehat{F}(n)}=\widehat{\bar{F}}(-n)$ for all $n \in \mathbb{Z}$, we conclude that $\widehat{F}(n)=0$ for all $n \in \mathbb{Z} \backslash\{0\}$, that is, $F$ is constant. Let us denote this constant by $\lambda$. Then (3.9) implies that

$$
\begin{equation*}
u_{+}=\bar{\lambda} \Psi=|\lambda|^{2} G . \tag{3.10}
\end{equation*}
$$

Since $\left\|u_{+}\right\|_{X^{\prime}}=1$, one gets $G \in X^{\prime}$ and $|\lambda|^{2}=\|G\|_{X^{\prime}}^{-1}$. It now follows from (3.7) and (3.10) that

$$
\|G\|_{X}=\frac{|\lambda|^{2}}{2 \pi} \int_{-\pi}^{\pi} G\left(e^{i \theta}\right) \overline{G\left(e^{i \theta}\right)} d \theta=\frac{1}{\|G\|_{X^{\prime}}} \frac{1}{2 \pi} \int_{-\pi}^{\pi}\left|G\left(e^{i \theta}\right)\right|^{2} d \theta=\frac{\|G\|_{L^{2}}^{2}}{\|G\|_{X^{\prime}}}
$$

Thus

$$
\|G\|_{X}\|G\|_{X^{\prime}}=\|G\|_{L^{2}}^{2}
$$

So, for every $g \in X_{b}$ such that $|g| \geq \varepsilon$ a.e. on $\mathbb{T}$, one has $g \in X^{\prime}, g \in L^{2}$, and

$$
\begin{equation*}
\|g\|_{X}\|g\|_{X^{\prime}}=\|g\|_{L^{2}}^{2} \tag{3.11}
\end{equation*}
$$

Take any $g \in X_{b}$ and set

$$
\Omega_{n}:=\left\{\zeta \in \mathbb{T}:|g(\zeta)| \geq \frac{1}{n}\right\}, \quad h_{n}:=g \chi_{\Omega_{n}}+\frac{1}{n} \chi_{\mathbb{T} \backslash \Omega_{n}}, \quad n \in \mathbb{N} .
$$

Then

$$
\left\|g-h_{n}\right\|_{X} \leq \frac{2}{n}\left\|\chi_{\mathbb{T} \backslash \Omega_{n}}\right\|_{X} \leq \frac{2}{n}\left\|\mathbf{e}_{0}\right\|_{X}, \quad n \in \mathbb{N}
$$

Hence $\left\{h_{n}\right\}$ converges to $g$ in $X$ as $n \rightarrow \infty$.
On the other hand, it is not difficult to see that for all $m, n \in \mathbb{N}$,

$$
\left|h_{m}-h_{n}\right| \leq \frac{2}{\min \{m, n\}} \quad \text { a.e. on } \quad \mathbb{T} \text {. }
$$

This implies that $\left\{h_{n}\right\}$ is a Cauchy sequence in $Y$, where $Y$ stands for $X^{\prime}$ or $L^{2}$. Since $X$ and $Y$ are continuously embedded into $L^{1}$, one concludes that $\left\{h_{n}\right\}$ converges to $g$ in $Y$. So,

$$
\begin{equation*}
X_{b} \subseteq X^{\prime} \cap L^{2} \tag{3.12}
\end{equation*}
$$

It follows from the above that (3.11) holds with $h_{n}$ in place of $g$. Passing to the limit as $n \rightarrow \infty$, one gets (3.11) for every $g \in X_{b}$.

Take now any $v \in L^{2}$. Let

$$
\chi_{0}:=\chi_{\{\zeta \in \mathbb{T}:|v(\zeta)|<1\}}, \quad \chi_{1}:=\chi_{\{\zeta \in \mathbb{T}:|v(\zeta)| \geq 1\}} .
$$

Since $\left|v \chi_{0}\right| \leq 1$ a.e. on $\mathbb{T}$, we have $v \chi_{0} \in X$. Let us show that $v \chi_{1} \in X$. If $\chi_{1}=0$ a.e. on $\mathbb{T}$, then there is nothing to prove. Otherwise, put

$$
v_{n}:=\chi_{1} \min \{|v|, n\}, \quad n \in \mathbb{N}
$$

Then $v_{n} \in L^{\infty} \subset X$ and

$$
\left\|v_{n}\right\|_{X}\left\|\chi_{1}\right\|_{X^{\prime}} \leq\left\|v_{n}\right\|_{X}\left\|v_{n}\right\|_{X^{\prime}}=\left\|v_{n}\right\|_{L^{2}}^{2}
$$

(see (3.11)). Since $v_{n} \uparrow|v| \chi_{1}$ as $n \rightarrow \infty$, it follows from [3, Ch. 1, Lemma 1.5] that $v \chi_{1} \in X$ and

$$
\left\|v \chi_{1}\right\|_{X} \leq \frac{\left\|v \chi_{1}\right\|_{L^{2}}^{2}}{\left\|\chi_{1}\right\|_{X^{\prime}}}
$$

Therefore $v=v \chi_{0}+v \chi_{1} \in X$. So, taking into account (3.12), we conclude that

$$
X_{b} \subseteq L^{2} \subseteq X
$$

Since $X_{b}, X$ and $L^{2}$ are continuously embedded into $L^{1}$, it follows from the closed graph theorem that the embeddings $X_{b} \subseteq L^{2} \subseteq X$ are also continuous and

$$
\begin{equation*}
C_{1}\|g\|_{X} \leq\|g\|_{L^{2}} \leq C_{2}\|g\|_{X} \quad \text { for all } \quad g \in X_{b} \tag{3.13}
\end{equation*}
$$

holds with some constants $C_{1}, C_{2} \in(0, \infty)$.
Finally, take any $g \in X$ and set $g_{n}:=\min \{|g|, n\}$. Then $g_{n} \in L^{\infty}$ and $g_{n} \uparrow|g|$ as $n \rightarrow \infty$. Hence, by Fatou's lemma (see [3, Ch. 1, Lemma 5.1]), $\left\|g_{n}\right\|_{X} \uparrow\|g\|_{X}<\infty$, and it follows from (3.13) that $\left\|g_{n}\right\|_{L^{2}} \uparrow \varrho<\infty$ for some constant $\varrho \leq C_{2}\|g\|_{X}$. So, $|g| \in L^{2}$, i.e. $g \in L^{2}$ for every $g \in X$, i.e. $X \subseteq L^{2}$. We conclude that $X=L^{2}$ and (3.11), (3.13) hold for all $g \in X$ (cf. [3, Ch. 1, Corollary 1.9]). It is now left to apply Lemma 3.1.

## 4. Necessary condition for the boundedness of the Hilbert transform from $C$ to a Banach function space $\boldsymbol{X}$

### 4.1. Operators of weak type from $C$ to a Banach function space $X$

Let $\mathcal{M}_{0}$ denote the subset of all almost everywhere finite functions in $\mathcal{M}$. It is well known (see, e.g., [14, Theorems 29.4.3 and 29.4.6] or [3, Ch. 1, Exercise 1]) that $\mathcal{M}_{0}$ can be equipped with a metric $d$ so that $\left(\mathcal{M}_{0}, d\right)$ is a complete linear metric space and the convergence in this metric is equivalent to the convergence in measure. Let $X$ be a Banach function spaces over the unit circle. We say that a linear operator $A: C \rightarrow \mathcal{M}_{0}$ is of weak type $(C, X)$ if there exists a constant $L>0$ such that for all $\lambda>0$ and all $f \in C$,

$$
\begin{equation*}
\left\|\chi_{\{\zeta \in \mathbb{T}:|(A f)(\zeta)|>\lambda\}}\right\|_{X} \leq L \frac{\|f\|_{C}}{\lambda} \tag{4.1}
\end{equation*}
$$

We denote the infimum of the constants $L$ satisfying (4.1) by $\|A\|_{C \rightarrow X}^{\text {weak }}$ and the set of all operators of weak type $(C, X)$ by $\mathcal{W}(C, X)$.

The proof of the following lemma is the same as that of [23, Lemma 3.1].
Lemma 4.1. Let $X$ be a Banach function space over the unit circle $\mathbb{T}$. If $A: C \rightarrow X$ is bounded, then $A \in \mathcal{W}(C, X)$ and $\|A\|_{C \rightarrow X}^{\text {weak }} \leq\|A\|_{C \rightarrow X}$.

### 4.2. Mapping of a finite family of separated arcs to a single arc

We will say that two open arcs in $\mathbb{T}$ are separated if the distance between them is positive, i.e. if they are disjoint and do not have common endpoints.

Theorem 4.2. If $E \subset \mathbb{T}$ is a finite union of pairwise separated open arcs,

$$
E=\bigcup_{k=1}^{n}\left(e^{i a_{k}}, e^{i b_{k}}\right) \neq \emptyset
$$

and $\ell \subset \mathbb{T}$ is an open arc such that $m(\ell)=m(E)$, then there exists a finite Blaschke product $u$ satisfying $u(0)=0$ and such that $u^{-1}(\ell)=E$.

Proof. The proof can easily be extracted from the proof of [7, Theorem 7.2] (note that the published version [8] of [7] contains a stronger variant of Theorem 7.2 equipped with a different proof that came from [31, Lemma 5.1]). We provide details here for the sake of completeness as a detailed proof of (4.4) (see below) was omitted in [7].

Take $\omega \in \mathbb{T} \backslash \operatorname{clos} E$ and consider

$$
\varphi(z):=i \frac{\omega+z}{\omega-z}
$$

This is a conformal homeomorphism of the unit disk $\mathbb{D}$ onto the upper half-plane $\mathbb{C}_{+}:=$ $\{\zeta \in \mathbb{C}: \operatorname{Im} \zeta>0\}$ and a diffeomorphism from $\mathbb{T} \backslash\{\omega\}$ onto $\mathbb{R}$. Let

$$
K(\zeta):=\prod_{k=1}^{n} \frac{\left|i-\varphi\left(e^{i a_{k}}\right)\right|}{\left|i-\varphi\left(e^{i b_{k}}\right)\right|} \cdot \frac{\zeta-\varphi\left(e^{i b_{k}}\right)}{\zeta-\varphi\left(e^{i a_{k}}\right)} .
$$

Then $K$ maps $\mathbb{C}_{+}$into itself, $\mathbb{R} \backslash\left\{\varphi\left(e^{i a_{k}}\right)\right\}_{k=1}^{n}$ into $\mathbb{R}$, and

$$
K^{-1}((-\infty, 0))=\bigcup_{k=1}^{n}\left(\varphi\left(e^{i a_{k}}\right), \varphi\left(e^{i b_{k}}\right)\right)
$$

(see [4, Proposition 2.1, Part (3)]).
If $\ell$ is an arc such that $m(\ell)=m(E)$, then there exists $a \in \mathbb{R}$ such that $\ell=$ $\left(e^{i a}, e^{i(a+2 \pi m(E))}\right)$. Let

$$
\psi(v):=e^{i a} \frac{v-e^{i \pi m(E)}}{v-e^{-i \pi m(E)}}
$$

Then $\psi$ is a conformal homeomorphism of $\mathbb{C}_{+}$onto $\mathbb{D}$ and a diffeomorphism from $\mathbb{R}$ onto $\mathbb{T} \backslash\left\{e^{i a}\right\}$ (see, e.g., [2, Theorem 13.16]). It is easy to see that $\psi^{-1}(\ell)=(-\infty, 0)$.

Let

$$
u:=\psi \circ K \circ \varphi
$$

Clearly, $u$ is a rational function. It is analytic in $\mathbb{D}$ and maps $\mathbb{D}$ into itself and $\mathbb{T} \backslash$ $\left(\{\omega\} \cup\left\{e^{i a_{k}}\right\}_{k=1}^{n}\right)$ into $\mathbb{T} \backslash\left\{e^{i a}\right\}$. The latter implies that $u$ does not have poles in $\mathbb{T}$ and hence is also analytic in a neighbourhood of $\mathbb{T}$. Therefore

$$
\lim _{|z| \rightarrow 1^{-}}|u(z)|=1
$$

It follows from [13, Theorem 3.5.2] (see also [13, Lemma 13.1.4] and [15, Ch. II, Sect. 6]) that $u$ is a finite Blaschke product.

We have

$$
\begin{aligned}
u^{-1}(\ell) & =\varphi^{-1}\left(K^{-1}\left(\psi^{-1}(\ell)\right)\right)=\varphi^{-1}\left(K^{-1}((-\infty, 0))\right) \\
& =\varphi^{-1}\left(\bigcup_{k=1}^{n}\left(\varphi\left(e^{i a_{k}}\right), \varphi\left(e^{i b_{k}}\right)\right)\right)=\bigcup_{k=1}^{n}\left(e^{i a_{k}}, e^{i b_{k}}\right)=E
\end{aligned}
$$

It is now left to show that $u(0)=0$.
Since $\varphi$ is a fractional linear transformation, it preserves the cross-ratio of any four points (see, e.g., [2, Theorem 13.23]). So,

$$
\frac{\left(\varphi(z)-\varphi\left(e^{i a_{k}}\right)\right)\left(\varphi(0)-\varphi\left(e^{i b_{k}}\right)\right)}{\left(\varphi(z)-\varphi\left(e^{i b_{k}}\right)\right)\left(\varphi(0)-\varphi\left(e^{i a_{k}}\right)\right)}=\frac{\left(z-e^{i a_{k}}\right) e^{i b_{k}}}{\left(z-e^{i b_{k}}\right) e^{i a_{k}}}, \quad k=1, \ldots, n
$$

Taking the limits as $z \rightarrow \omega$, we get

$$
\begin{equation*}
\frac{i-\varphi\left(e^{i b_{k}}\right)}{i-\varphi\left(e^{i a_{k}}\right)}=\frac{\left(\omega-e^{i a_{k}}\right) e^{i b_{k}}}{\left(\omega-e^{i b_{k}}\right) e^{i a_{k}}}, \quad k=1, \ldots, n \tag{4.2}
\end{equation*}
$$

By the inscribed angle theorem, the angle at $\omega$ subtended by the $\operatorname{arc}\left(e^{i a_{k}}, e^{i b_{k}}\right)$ is equal to $\left(b_{k}-a_{k}\right) / 2$. Hence

$$
\begin{equation*}
\frac{e^{i b_{k}}-\omega}{\left|e^{i b_{k}}-\omega\right|}\left(\frac{e^{i a_{k}}-\omega}{\left|e^{i a_{k}}-\omega\right|}\right)^{-1}=e^{i\left(b_{k}-a_{k}\right) / 2}, \quad k=1, \ldots, n \tag{4.3}
\end{equation*}
$$

Taking into account (4.2)-(4.3), we get

$$
\begin{aligned}
K(i) & =\prod_{k=1}^{n} \frac{\left|i-\varphi\left(e^{i a_{k}}\right)\right|}{\left|i-\varphi\left(e^{i b_{k}}\right)\right|} \cdot \frac{i-\varphi\left(e^{i b_{k}}\right)}{i-\varphi\left(e^{i a_{k}}\right)}=\prod_{k=1}^{n}\left|\frac{\omega-e^{i b_{k}}}{\omega-e^{i a_{k}}}\right| \cdot \frac{\omega-e^{i a_{k}}}{\omega-e^{i b_{k}}} e^{i\left(b_{k}-a_{k}\right)} \\
& =\prod_{k=1}^{n} \frac{e^{i a_{k}}-\omega}{\left|e^{i a_{k}}-\omega\right|}\left(\frac{e^{i b_{k}}-\omega}{\left|e^{i b_{k}}-\omega\right|}\right)^{-1} e^{i\left(b_{k}-a_{k}\right)}=\prod_{k=1}^{n} e^{-i\left(b_{k}-a_{k}\right) / 2} e^{i\left(b_{k}-a_{k}\right)}
\end{aligned}
$$

$$
\begin{equation*}
=\prod_{k=1}^{n} e^{i\left(b_{k}-a_{k}\right) / 2}=\exp \left(i \sum_{k=1}^{n}\left(b_{k}-a_{k}\right) / 2\right)=e^{i \pi m(E)} . \tag{4.4}
\end{equation*}
$$

Hence

$$
u(0)=\psi(K(\varphi(0)))=\psi(K(i))=\psi\left(e^{i \pi m(E)}\right)=0
$$

which completes the proof.

### 4.3. Estimates for the Hilbert transform of a piecewise linear bump function

In this subsection, we prove a lower estimate for the Hilbert transform of a piecewise linear bump function, which will play an important role in what follows.

Lemma 4.3. Let $0<\beta<\frac{\pi}{2}, 0<\varepsilon<\min \left\{\beta, \frac{\pi}{2}-\beta\right\}$, and

$$
g\left(e^{i \theta}\right):= \begin{cases}(\theta+\pi) / \varepsilon, & -\pi \leq \theta \leq-\pi+\varepsilon  \tag{4.5}\\ 1, & -\pi+\varepsilon<\theta<-\beta \\ -(\theta+\beta-\varepsilon) / \varepsilon, & -\beta \leq \theta \leq-\beta+\varepsilon \\ 0, & -\beta+\varepsilon<\theta \leq \pi\end{cases}
$$

Then

$$
\begin{equation*}
\left|(H g)\left(e^{i \eta}\right)\right|>\frac{1}{\pi}\left|\log \left(\sqrt{2} \sin \frac{\beta+\varepsilon}{2}\right)\right|-\frac{\varepsilon}{2 \pi} \quad \text { for all } \quad \eta \in[\pi-\beta, \pi] \tag{4.6}
\end{equation*}
$$

Proof. Take any $\eta \in[\pi-\beta, \pi]$. Since $\cot \frac{\eta-\theta}{2} \leq 0$ for $\theta \in[-\pi,-\beta]$, we have

$$
\begin{align*}
(H g)\left(e^{i \eta}\right) & =\frac{1}{2 \pi} \int_{-\pi}^{\pi} g\left(e^{i \theta}\right) \cot \frac{\eta-\theta}{2} d \theta \\
& =\frac{1}{2 \pi}\left(\int_{-\pi}^{-\pi+\varepsilon}+\int_{-\pi+\varepsilon}^{-\beta}+\int_{-\beta}^{-\beta+\varepsilon} g\left(e^{i \theta}\right) \cot \frac{\eta-\theta}{2} d \theta\right. \\
& \leq \frac{1}{2 \pi} \int_{-\pi+\varepsilon}^{-\beta} \cot \frac{\eta-\theta}{2} d \theta+\frac{1}{2 \pi} \int_{-\beta}^{-\beta+\varepsilon}\left|\cot \frac{\eta-\theta}{2}\right| d \theta . \tag{4.7}
\end{align*}
$$

An easy calculation shows that

$$
\begin{equation*}
\frac{1}{2 \pi} \int_{-\pi+\varepsilon}^{-\beta} \cot \frac{\eta-\theta}{2} d \theta=\frac{1}{\pi} \log \sin \frac{\eta+\pi-\varepsilon}{2}-\frac{1}{\pi} \log \sin \frac{\eta+\beta}{2} \tag{4.8}
\end{equation*}
$$

Since

$$
\frac{\pi}{2}<\pi-\frac{\beta+\varepsilon}{2} \leq \frac{\eta+\pi-\varepsilon}{2} \leq \pi-\frac{\varepsilon}{2}<\pi
$$

for $\eta \in[\pi-\beta, \pi]$, we have

$$
\begin{equation*}
\frac{1}{\pi} \log \sin \frac{\eta+\pi-\varepsilon}{2} \leq \frac{1}{\pi} \log \sin \left(\pi-\frac{\beta+\varepsilon}{2}\right)=\frac{1}{\pi} \log \sin \frac{\beta+\varepsilon}{2} \tag{4.9}
\end{equation*}
$$

Similarly, the inequalities

$$
\frac{\pi}{2} \leq \frac{\eta+\beta}{2} \leq \frac{\pi+\beta}{2}<\pi
$$

for $\eta \in[\pi-\beta, \pi]$, imply that

$$
\begin{equation*}
\frac{1}{\pi} \log \sin \frac{\eta+\beta}{2} \geq \frac{1}{\pi} \log \sin \frac{\pi+\beta}{2}=\frac{1}{\pi} \log \cos \frac{\beta}{2} . \tag{4.10}
\end{equation*}
$$

Taking into account

$$
0<\frac{\beta}{2}<\frac{\beta+\varepsilon}{2}<\frac{\pi}{4}
$$

we see that

$$
\sin \frac{\beta+\varepsilon}{2}<\frac{1}{\sqrt{2}}<\cos \frac{\beta}{2}
$$

Hence

$$
\begin{equation*}
-\frac{1}{\pi} \log \cos \frac{\beta}{2}<\frac{1}{\pi} \log \sqrt{2}<-\frac{1}{\pi} \log \sin \frac{\beta+\varepsilon}{2} . \tag{4.11}
\end{equation*}
$$

Combining (4.8)-(4.11), we arrive at

$$
\begin{align*}
\frac{1}{2 \pi} \int_{-\pi+\varepsilon}^{-\beta} \cot \frac{\eta-\theta}{2} d \theta & <\frac{1}{\pi} \log \sin \frac{\beta+\varepsilon}{2}+\frac{1}{\pi} \log \sqrt{2} \\
& =\frac{1}{\pi} \log \left(\sqrt{2} \sin \frac{\beta+\varepsilon}{2}\right) \\
& =-\frac{1}{\pi}\left|\log \left(\sqrt{2} \sin \frac{\beta+\varepsilon}{2}\right)\right| \tag{4.12}
\end{align*}
$$

Since

$$
\frac{\pi-\varepsilon}{2} \leq \frac{\eta+\beta-\varepsilon}{2} \leq \frac{\eta-\theta}{2} \leq \frac{\eta+\beta}{2} \leq \frac{\pi+\beta}{2}
$$

for $\eta \in[\pi-\beta, \pi]$ and $\theta \in[-\beta,-\beta+\varepsilon]$, the function $\cot \varphi$ is decreasing for $\varphi \in\left[\frac{\pi-\varepsilon}{2}, \frac{\pi+\beta}{2}\right]$, and

$$
0<\cot \frac{\pi-\varepsilon}{2}=-\cot \frac{\pi+\varepsilon}{2}<-\cot \frac{\pi+\beta}{2}=\tan \frac{\beta}{2}
$$

we have

$$
\begin{align*}
\frac{1}{2 \pi} \int_{-\beta}^{-\beta+\varepsilon}\left|\cot \frac{\eta-\theta}{2}\right| d \theta & \leq \frac{\varepsilon}{2 \pi} \max _{\varphi \in\left[\frac{\pi-\varepsilon}{2}, \frac{\pi+\beta}{2}\right]}|\cot \varphi| \\
& =\frac{\varepsilon}{2 \pi} \tan \frac{\beta}{2} \leq \frac{\varepsilon}{2 \pi} \tan \frac{\pi}{4}=\frac{\varepsilon}{2 \pi} \tag{4.13}
\end{align*}
$$

It follows from (4.7), (4.12), and (4.13) that for $\eta \in[\pi-\beta, \pi]$,

$$
-\left|(H g)\left(e^{i \eta}\right)\right| \leq(H g)\left(e^{i \eta}\right)<-\frac{1}{\pi}\left|\log \left(\sqrt{2} \sin \frac{\beta+\varepsilon}{2}\right)\right|+\frac{\varepsilon}{2 \pi}
$$

which immediately implies (4.6).

### 4.4. Necessary conditions for the Hilbert transform to be of weak type $(C, X)$

In this subsection, we show that if the Hilbert transform is of weak type $(C, X)$ for some Banach function space $X$, then $X_{a}=X_{b}$.

Let $E$ be a union of pairwise disjoint arcs of small measure. Then

$$
\begin{equation*}
F(m(E) ; \varepsilon):=\frac{1}{\pi}\left|\log \left(\sqrt{2} \sin \left(\pi m(E)+\frac{\varepsilon}{2}\right)\right)\right|-\frac{\varepsilon}{2 \pi} \tag{4.14}
\end{equation*}
$$

is large whenever $\varepsilon>0$ is small. We start by constructing a continuous real-valued function $f$ depending on $\varepsilon$ such that $|f| \leq 1$ while the modulus of the Hilbert transform of $f$ exceeds $F(m(E) ; \varepsilon)$ on the set $E$. This function is the composition of the piecewise linear bump function from Lemma 4.3 and the finite Blaschke product from Theorem 4.2.

Lemma 4.4. Let $E \subset \mathbb{T}$ be a finite union of pairwise disjoint open arcs such that $0<m(E)<\frac{1}{4}$. Then for every positive $\varepsilon<2 \pi \min \left\{m(E), \frac{1}{4}-m(E)\right\}$ there exists a continuous function $f: \mathbb{T} \rightarrow \mathbb{R}$ such that $|f| \leq 1$ and

$$
\begin{equation*}
\left|(H f)\left(e^{i \eta}\right)\right|>\frac{1}{\pi}\left|\log \left(\sqrt{2} \sin \left(\pi m(E)+\frac{\varepsilon}{2}\right)\right)\right|-\frac{\varepsilon}{2 \pi} \text { for all } e^{i \eta} \in E . \tag{4.15}
\end{equation*}
$$

Proof. If the pairwise disjoint open arcs constituting $E$ are pairwise separated, set $E_{0}:=$ $E$. Otherwise, let $E_{0}$ be the union of $E$ with the set of common endpoints of the nonseparated arcs in the family constituting $E$. In this case, $E_{0}$ is obtained from $E$ by
merging adjacent open arcs into bigger ones and reducing the total number of arcs. Either way, $E_{0}$ is a finite union of pairwise separated open arcs, $E \subseteq E_{0}$, and $m\left(E_{0}\right)=m(E)$.

Let $g$ be defined by (4.5) with $\beta=2 \pi m(E), u$ be the finite Blaschke product from Theorem 4.2 with $\ell=\left\{e^{i \theta} \in \mathbb{T}: \theta \in(\pi-2 \pi m(E), \pi)\right\}$ and with $E_{0}$ in place of $E$. Consider $f:=g \circ u$. Since $g$ and $u$ are continuous on $\mathbb{T}$, so is $f$. Since $H f=H(g \circ u)=$ $(\mathrm{Hg}) \circ u$ in view of Lemma 2.4, it follows from Lemma 4.3 that

$$
\left|(H f)\left(e^{i \eta}\right)\right|=\left|(H g)\left(u\left(e^{i \eta}\right)\right)\right|>\frac{1}{\pi}\left|\log \left(\sqrt{2} \sin \left(\pi m\left(E_{0}\right)+\frac{\varepsilon}{2}\right)\right)\right|-\frac{\varepsilon}{2 \pi}
$$

for all $e^{i \eta}$ such that $u\left(e^{i \eta}\right) \in \ell$, i.e. for all $e^{i \eta} \in u^{-1}(\ell)=E_{0}$. This immediately implies (4.15).

Corollary 4.5. Let $E \subset \mathbb{T}$ be a finite union of pairwise disjoint open arcs such that $0<m(E)<\frac{1}{4}$. Then there exists a continuous function $f: \mathbb{T} \rightarrow \mathbb{R}$ such that $|f| \leq 1$ and

$$
\begin{equation*}
\left|(H f)\left(e^{i \eta}\right)\right|>\frac{1}{2 \pi}|\log (\sqrt{2} \sin (\pi m(E)))| \quad \text { for all } \quad e^{i \eta} \in E \tag{4.16}
\end{equation*}
$$

Proof. Let $F(m(E), \varepsilon)$ be defined by (4.14). Since it is continuous in $\varepsilon$, there exists $\varepsilon>0$ such that $F(m(E), \varepsilon)-F(m(E), 0)>F(m(E), 0) / 2$, whence

$$
F(m(E), \varepsilon)>F(m(E), 0) / 2
$$

By Lemma 4.4, there exists a continuous real-valued function $f$ such that (4.15) holds. Combining (4.15) with the above inequality, we arrive at (4.16).

Next we show that if $E$ is a finite union of pairwise disjoint open arcs of small measure and $H \in \mathcal{W}(C, X)$, then $\left\|\chi_{E}\right\|_{X}=O(1 / F(m(E), 0))$.

Lemma 4.6. Let $X$ be a Banach function space over the unit circle $\mathbb{T}$. If the Hilbert transform $H$ is of weak type $(C, X)$, then for every finite union $E$ of pairwise disjoint open arcs such that $0<m(E)<\frac{1}{4}$, one has

$$
\begin{equation*}
\left\|\chi_{E}\right\|_{X} \leq \frac{2 \pi\|H\|_{C \rightarrow X}^{\text {weak }}}{|\log (\sqrt{2} \sin (\pi m(E)))|} \tag{4.17}
\end{equation*}
$$

Proof. Let

$$
\lambda=\frac{1}{2 \pi}|\log (\sqrt{2} \sin (\pi m(E)))|
$$

By Corollary 4.5 , there exists a function $f \in C$ such that $|f| \leq 1$ and

$$
\left.\chi_{E}(\tau) \leq \chi_{\{\zeta \in \mathbb{T}}:|(H f)(\zeta)|>\lambda\right\}(\tau), \quad \tau \in \mathbb{T}
$$

Therefore, by the lattice property, taking into account that $H \in \mathcal{W}(C, X)$, we obtain

$$
\begin{aligned}
\left\|\chi_{E}\right\|_{X} & \leq\left\|\chi_{\{\zeta \in \mathbb{T}:|(H f)(\zeta)|>\lambda\}}\right\|_{X} \leq \frac{1}{\lambda}\|H\|_{C \rightarrow X}^{\text {weak }}\|f\|_{C} \\
& \leq \frac{2 \pi\|H\|_{C \rightarrow X}^{\text {weak }}}{|\log (\sqrt{2} \sin (\pi m(E)))|}
\end{aligned}
$$

which completes the proof.

Now we are in a position to prove the main result of this section.

Theorem 4.7. Let $X$ be a Banach function space over the unit circle $\mathbb{T}$. If the Hilbert transform $H$ is of weak type $(C, X)$, then $X_{a}=X_{b}$.

Proof. Consider a sequence $\left\{\gamma_{j}\right\}_{j \in \mathbb{N}}$ of measurable subsets of $\mathbb{T}$ such that $\chi_{\gamma_{j}} \rightarrow 0$ a.e. on $\mathbb{T}$ as $j \rightarrow \infty$. By the dominated convergence theorem,

$$
m\left(\gamma_{j}\right)=\int_{\mathbb{T}} \chi_{\gamma_{j}}(\tau) d m(\tau) \rightarrow 0 \quad \text { as } \quad j \rightarrow \infty
$$

Without loss of generality, one can assume that $0<m\left(\gamma_{j}\right)<\frac{1}{8}$ for all $j \in \mathbb{N}$. For every $j \in \mathbb{N}$, there exists an open set $\mathcal{E}_{j}$ such that $\gamma_{j} \subseteq \mathcal{E}_{j}$ and $m\left(\mathcal{E}_{j}\right) \leq 2 m\left(\gamma_{j}\right)$. Each $\mathcal{E}_{j}$ is the union of an at most countable family of pairwise disjoint open arcs:

$$
\mathcal{E}_{j}=\bigcup_{k=1}^{N_{j}} \ell_{j, k}, \quad N_{j} \in \mathbb{N} \cup\{\infty\}
$$

If $N_{j}$ is finite, set $E_{j}:=\mathcal{E}_{j}$. Otherwise, let $\mathcal{E}_{j}^{n}=\bigcup_{k=1}^{n} \ell_{j, k}$. Since $\chi_{\mathcal{E}_{j}^{n}} \uparrow \chi_{\mathcal{E}_{j}}$ a.e. as $n \rightarrow \infty$, it follows from the Fatou property (A3) that

$$
\left\|\chi_{\mathcal{E}_{j}^{n}}\right\|_{X} \uparrow\left\|\chi_{\mathcal{E}_{j}}\right\|_{X} \quad \text { as } \quad n \rightarrow \infty
$$

Then there exists $n_{j} \in \mathbb{N}$ such that

$$
\left\|\chi_{\mathcal{E}_{j}^{n_{j}}}\right\|_{X} \geq \frac{1}{2}\left\|\chi_{\mathcal{E}_{j}}\right\|_{X}
$$

Set $E_{j}:=\mathcal{E}_{j}^{n_{j}}$. Then $E_{j}$ is a finite union of pairwise disjoint open arcs,

$$
\frac{1}{2} m\left(E_{j}\right) \leq \frac{1}{2} m\left(\mathcal{E}_{j}\right) \leq m\left(\gamma_{j}\right)<\frac{1}{8}, \quad\left\|\chi_{E_{j}}\right\|_{X} \geq \frac{1}{2}\left\|\chi_{\mathcal{E}_{j}}\right\|_{X} \geq \frac{1}{2}\left\|\chi_{\gamma_{j}}\right\|_{X}
$$

By Lemma 4.6, for every $j \in \mathbb{N}$,

$$
\left\|\chi_{\gamma_{j}}\right\|_{X} \leq 2\left\|\chi_{E_{j}}\right\|_{X} \leq \frac{4 \pi\|H\|_{C \rightarrow X}^{\text {weak }}}{\left|\log \left(\sqrt{2} \sin \left(\pi m\left(E_{j}\right)\right)\right)\right|} \leq \frac{4 \pi\|H\|_{C \rightarrow X}^{\text {weak }}}{\left|\log \left(\sqrt{2} \sin \left(2 \pi m\left(\gamma_{j}\right)\right)\right)\right|}
$$

Since $m\left(\gamma_{j}\right) \rightarrow 0$ as $j \rightarrow \infty$, the above estimate implies that $\left\|\chi_{\gamma_{j}}\right\|_{X} \rightarrow 0$ as $j \rightarrow \infty$. Thus the constant function 1 has absolutely continuous norm. Then it follows from [3, Ch. 1, Theorem 3.8] that for every measurable set $E \subset \mathbb{T}$, its characteristic function $\chi_{E}$ has absolutely continuous norm. Thus, by [3, Ch. 1, Theorem 3.13], $X_{a}=X_{b}$.

The above theorem and Lemma 4.1 immediately imply the following.

Corollary 4.8. Let $X$ be a Banach function space over the unit circle $\mathbb{T}$. If the Hilbert transform $H$ is bounded from the space of continuous functions $C$ to a Banach function space $X$, then $X_{a}=X_{b}$.

## 5. Proof of the second main result

5.1. Necessary condition for the boundedness of the Riesz projection from $C$ to a Banach function space $X$

We start this section by rephrasing Corollary 4.8 in terms of the Riesz projection. It improves [23, Theorem 3.7] and [25, Theorem 3.3].

Theorem 5.1. If the Riesz projection $P$ is bounded from the space of continuous functions $C$ to a Banach function space $X$, then $X_{a}=X_{b}$.

Proof. If $f \in C \subset L^{1}$, then

$$
\begin{equation*}
P f:=\frac{1}{2}(f+i H f)+\frac{1}{2} \widehat{f}(0) \tag{5.1}
\end{equation*}
$$

(cf. [15, p. 104], [6, Section 1.43] and also [23, formula (1.3)]). Since $C$ is continuously embedded into $L^{1}$, the functional $f \mapsto \widehat{f}(0)$ is continuous on the space $C$. Then it follows from (5.1) that $P: C \rightarrow X$ is bounded if and only if $H: C \rightarrow X$ is bounded. It follows from this observation and Corollary 4.8 that $X_{a}=X_{b}$.

### 5.2. A relation between the backward shift and the Riesz projection

The next lemma relates the backward shift operator with the Riesz projection.

Lemma 5.2. If $f \in H^{1}$, then

$$
\begin{equation*}
S f=P\left(\mathbf{e}_{-1} f\right) \tag{5.2}
\end{equation*}
$$

Proof. Lemma 2.6 implies that $P f=f$. Hence

$$
\begin{aligned}
\left(P\left(\mathbf{e}_{-1} f\right)\right)(t) & =\frac{\left(\mathbf{e}_{-1} f\right)(t)+\left(\mathcal{C}\left(\mathbf{e}_{-1} f\right)\right)(t)}{2} \\
& \left.=\mathbf{e}_{-1}(t) \frac{f(t)+(\mathcal{C} f)(t)}{2}+\frac{1}{2}\left(\mathcal{C}\left(\mathbf{e}_{-1} f\right)\right)(t)-\mathbf{e}_{-1}(t)(\mathcal{C} f)(t)\right) \\
& =\mathbf{e}_{-1}(t)(P f)(t)+\frac{1}{2 \pi i} \text { p.v. } \int_{\mathbb{T}}\left(\frac{1}{\tau}-\frac{1}{t}\right) \frac{f(\tau)}{\tau-t} d \tau \\
& =\mathbf{e}_{-1}(t) f(t)-\frac{1}{2 \pi i} \int_{\mathbb{T}} \frac{f(\tau)}{\tau t} d \tau=\mathbf{e}_{-1}(t) f(t)-\frac{\mathbf{e}_{-1}(t)}{2 \pi} \int_{-\pi}^{\pi} f\left(e^{i \theta}\right) d \theta \\
& =\mathbf{e}_{-1}(t) f(t)-\mathbf{e}_{-1}(t) \widehat{f}(0)=\mathbf{e}_{-1}(t)(f(t)-\widehat{f}(0))=(S f)(t) .
\end{aligned}
$$

### 5.3. Necessary conditions for the backward shift operator to have norm one

In this subsection we point out a consequence of $\|S\|_{H[X] \rightarrow H[X]}=1$.
Lemma 5.3. If $X$ is a Banach function space and $\|S\|_{H[X] \rightarrow H[X]}=1$, then

$$
\begin{equation*}
\|P g\|_{X} \leq\|g\|_{X} \quad \text { for all } \quad g \in \mathcal{P} \tag{5.3}
\end{equation*}
$$

Proof. Let $Q:=I-P$. For any trigonometric polynomial

$$
g\left(e^{i \theta}\right)=\sum_{k=-M}^{N} g_{k} e^{i k \theta}
$$

where $M, N \in \mathbb{Z}_{+}$, one has

$$
P g=P\left(\mathbf{e}_{-1} \mathbf{e}_{1} g\right)=P\left(\mathbf{e}_{-1} P\left(\mathbf{e}_{1} g\right)\right)+P\left(\mathbf{e}_{-1} Q\left(\mathbf{e}_{1} g\right)\right)=P\left(\mathbf{e}_{-1} P\left(\mathbf{e}_{1} g\right)\right)
$$

Since $P\left(\mathbf{e}_{1} g\right) \in H[X] \cap \mathcal{P} \subset H^{1}$, it follows from Lemma 5.2 and the above equality that $P g=S P\left(\mathbf{e}_{1} g\right)$. Repeating the above argument $M$ times, we get

$$
P g=S^{M} P\left(\mathbf{e}_{M} g\right)
$$

Since $\mathbf{e}_{M} g \in H[X] \cap \mathcal{P}$, we have $P\left(\mathbf{e}_{M} g\right)=\mathbf{e}_{M} g$. Hence, taking into account that $\|S\|_{H[X] \rightarrow H[X]}=1$, we get

$$
\begin{aligned}
\|P g\|_{X} & =\|P g\|_{H[X]}=\left\|S^{M} P\left(\mathbf{e}_{M} g\right)\right\|_{H[X]} \leq\left\|S^{M}\right\|_{H[X] \rightarrow H[X]}\left\|P\left(\mathbf{e}_{M} g\right)\right\|_{H[X]} \\
& \leq\|S\|_{H[X] \rightarrow H[X]}^{M}\left\|P\left(\mathbf{e}_{M} g\right)\right\|_{H[X]}=\left\|\mathbf{e}_{M} g\right\|_{H[X]}=\|g\|_{X}
\end{aligned}
$$

which completes the proof of (5.3).

### 5.4. Proof of Theorem 1.2

Take $f \in H[X] \backslash\{0\}$ such that $\widehat{f}(0)=0$. Then $S f=\mathbf{e}_{-1} f$, which implies that $|S f|=|f|$ a.e. on $\mathbb{T}$. Hence $\|S f\|_{H[X]}=\|f\|_{H[X]}$. On the other hand, $f \in H[X] \subset H^{1}$. Then it follows from Lemma 2.6 that $P f=f$. So, one always has

$$
\begin{equation*}
\|S\|_{H[X] \rightarrow H[X]} \geq 1 \quad \text { and } \quad\|P\|_{X \rightarrow X} \geq 1 \tag{5.4}
\end{equation*}
$$

Necessity. Suppose $\|P\|_{X \rightarrow X}=1$. It follows from Lemma 5.2 that $S f=P\left(\mathbf{e}_{-1} f\right)$ for $f \in H[X] \subset H^{1}$. Hence

$$
\|S f\|_{H[X]}=\left\|P\left(\mathbf{e}_{-1} f\right)\right\|_{X} \leq\|P\|_{X \rightarrow X}\left\|\mathbf{e}_{-1} f\right\|_{X}=\|f\|_{X}
$$

Hence $\|S\|_{H[X] \rightarrow H[X]} \leq 1$. This inequality and the first inequality in (5.4) imply that $\|S\|_{H[X] \rightarrow H[X]}=1$.

Sufficiency. Suppose that $\|S\|_{H[X] \rightarrow H[X]}=1$. It follows from Lemma 5.3 and from Axioms (A2) and (A4) in the definition of a Banach function space that there exists $k>0$ such that

$$
\begin{equation*}
\|P g\|_{X} \leq\|g\|_{X} \leq k\|g\|_{C} \quad \text { for all } \quad g \in \mathcal{P} \tag{5.5}
\end{equation*}
$$

This inequality and the Weierstrass approximation theorem (see, e.g., corollary to [26, Ch. 1, Theorem 2.12]) imply that $P: C \rightarrow X$ is bounded. Then Theorem 5.1 implies that $X_{a}=X_{b}$. By Lemma 2.1, the set of trigonometric polynomials $\mathcal{P}$ is dense in $X_{b}$. Then the first inequality in (5.5) implies that

$$
\|P f\|_{X} \leq\|f\|_{X} \quad \text { for all } \quad f \in X_{b}
$$

Therefore $\|P\|_{X_{b} \rightarrow X} \leq 1$. Then Theorem 2.7 implies that

$$
\begin{equation*}
\|P\|_{X \rightarrow X}=\|P\|_{X_{b} \rightarrow X} \leq 1 \tag{5.6}
\end{equation*}
$$

Combining the second inequality in (5.4) with (5.6), we arrive at the equality $\|P\|_{X \rightarrow X}=$ 1.

Remark 5.4. Let

$$
\left(P_{0} f\right)(t):=f(t)-\widehat{f}(0), \quad t \in \mathbb{T}
$$

Since $\left|\left(P_{0} f\right)(t)\right|=|(S f)(t)|$, we have $\left\|P_{0}\right\|_{H[X] \rightarrow H[X]}=\|S\|_{H[X] \rightarrow H[X]}$. Hence it follows from Theorems 1.1 and 1.2 that if $\left\|P_{0}\right\|_{H[X] \rightarrow H[X]}=1$, then $X$ coincides with $L^{2}$ and (1.5) holds.

The same is true if $\left\|P_{0}\right\|_{X \rightarrow X}=1$, since $\left\|P_{0}\right\|_{H[X] \rightarrow H[X]}=\left\|P_{0}\right\|_{X \rightarrow X}=1$ in this case. Indeed, $P_{0} \mathbf{e}_{1}=\mathbf{e}_{1}$, whence $\left\|P_{0}\right\|_{H[X] \rightarrow H[X]} \geq 1$. On the other hand, $P_{0}: H[X] \rightarrow H[X]$ is the restriction of $P_{0}: X \rightarrow X$ to $H[X]$, and $1 \leq\left\|P_{0}\right\|_{H[X] \rightarrow H[X]} \leq\left\|P_{0}\right\|_{X \rightarrow X}=1$.

It is easy to see that $P_{0}: X \rightarrow X$ is a projection onto a subspace of codimension one, and it is instructive to compare the above results to the following ones.

Suppose that $X$ is a real separable Banach function space such that $\left\|P_{0}\right\|_{X \rightarrow X}=1$. It follows from [3, Ch. 1, Corollary 5.6] that $X=X_{a}$. Then, by [32, Theorem 2] (see also [20, Theorem 4.3]), there exists a positive measurable function $w$ such that

$$
\|g\|_{X}=\left(\int_{\mathbb{T}} g^{2}(t) w(t) d m(t)\right)^{1 / 2} \quad \text { for all } \quad g \in X
$$

In this case, $\left\|P_{0} g\right\|_{X} \leq\|g\|_{X}$ is equivalent to

$$
\begin{equation*}
(\widehat{g}(0))^{2} \int_{\mathbb{T}} w(t) d m(t) \leq 2 \widehat{g}(0) \int_{\mathbb{T}} g(t) w(t) d m(t) \tag{5.7}
\end{equation*}
$$

It is easy to see that if $w$ is non-constant, then there exists a simple function $g$ such that $\widehat{g}(0)>0$ and $\int_{\mathbb{T}} g(t) w(t) d m(t)=0$. For such a function, (5.7) cannot hold. So, $w$ has to be constant, which means that $X$ coincides with $L^{2}$ and (1.5) holds.

If $X$ is a real separable rearrangement-invariant Banach function space, and there exists a projection $Q: X \rightarrow X$ onto a subspace of finite codimension such that $\|Q\|_{X \rightarrow X}=1$, then $X$ is isometric to $L^{2}$ ([32, Theorem 4]), and hence $X$ coincides with $L^{2}$ and (1.5) holds (see [1, Theorem 1]).

## 6. The norm of the Riesz projection in terms of Toeplitz operators

Let $X$ be a Banach function space on which the Riesz projection $P$ is bounded. For $a \in L^{\infty}$, the Toeplitz operator $T(a)$ on $X$ is defined by

$$
T(a) f=P(a f), \quad f \in H[X] .
$$

It is easy to see that

$$
\|T(a)\|_{H[X] \rightarrow H[X]} \leq\|P\|_{X \rightarrow X}\|a\|_{L^{\infty}}
$$

Note that if $P$ is bounded on $X$, then in view of Lemma 5.2, the backward shift operator $S$ coincides with the Toeplitz operator $T\left(\mathbf{e}_{-1}\right)$ :

$$
S f=T\left(\mathbf{e}_{-1}\right) f, \quad f \in H[X] .
$$

Following [15, Ch. IX, Section 2], let

$$
C+H^{\infty}:=\left\{f+g: f \in C, g \in H^{\infty}\right\} .
$$

It is well known that $C+H^{\infty}$ is a closed subalgebra of $L^{\infty}$ generated by the set $H^{\infty} \cup$ $\left\{\mathbf{e}_{-1}\right\}$ (see, e.g., [15, Ch. IX, Theorem 2.2]).

The following theorem sharpens a part of [5, Theorem 7.7] and extends it from the setting of Lebesgue spaces $L^{p}$ to the setting of Banach function spaces $X$.

Theorem 6.1. Let $X$ be a Banach function space on which the Riesz projection is bounded and

$$
\begin{aligned}
c_{X} & :=\sup _{n \in \mathbb{N}}\left\|T\left(\mathbf{e}_{-n}\right)\right\|_{H[X] \rightarrow H[X]}, \\
s_{X} & :=\sup _{a \in\left(C+H^{\infty}\right) \backslash\{0\}} \frac{\|T(a)\|_{H[X] \rightarrow H[X]}}{\|a\|_{L^{\infty}}}, \\
\sigma_{X} & :=\sup _{a \in L^{\infty} \backslash\{0\}} \frac{\|T(a)\|_{H[X] \rightarrow H[X]}}{\|a\|_{L^{\infty}}} .
\end{aligned}
$$

Then

$$
c_{X}=s_{X}=\sigma_{X}=\|P\|_{X \rightarrow X}
$$

Proof. It is clear that

$$
c_{X} \leq s_{X} \leq \sigma_{X} \leq\|P\|_{X \rightarrow X}
$$

So, it is sufficient to show that

$$
\begin{equation*}
\|P\|_{X \rightarrow X} \leq c_{X} \tag{6.1}
\end{equation*}
$$

By Theorem 2.7, $\|P\|_{X \rightarrow X}=\|P\|_{X_{b} \rightarrow X}$. Hence for any $\varepsilon>0$ there exists $f \in X_{b}$ such that $\|f\|_{X}=1$ and

$$
\|P f\|_{X}>\|P\|_{X \rightarrow X}-\varepsilon
$$

Since $P: X \rightarrow X$ is bounded and $C$ is continuously embedded into $X$ (by Axioms (A2) and (A4) in the definition of a Banach function space), we see that $P: C \rightarrow X$ is bounded. Hence it follows from Theorem 5.1 that $X_{a}=X_{b}$. Then by Lemma 2.1, there exists a trigonometric polynomial

$$
g\left(e^{i \theta}\right)=\sum_{k=-M}^{N} g_{k} e^{i k \theta}
$$

such that $\|f-g\|_{X}<\varepsilon$. Then $\|g\|_{X}<1+\varepsilon$ and

$$
\|P g\|_{X} \geq\|P f\|_{X}-\|P(f-g)\|_{X}>\|P\|_{X \rightarrow X}-\varepsilon-\|P\|_{X \rightarrow X} \varepsilon
$$

Since $\mathbf{e}_{M} g \in H[X]$, one has

$$
T\left(\mathbf{e}_{-M}\right)\left(\mathbf{e}_{M} g\right)=P\left(\mathbf{e}_{-M} \mathbf{e}_{M} g\right)=P g
$$

Therefore

$$
\begin{aligned}
\left\|T\left(\mathbf{e}_{-M}\right)\right\|_{H[X] \rightarrow H[X]} & \geq \frac{\left\|T\left(\mathbf{e}_{-M}\right)\left(\mathbf{e}_{M} g\right)\right\|_{H[X]}}{\left\|\mathbf{e}_{M} g\right\|_{H[X]}}=\frac{\|P g\|_{H[X]}}{\|g\|_{X}}=\frac{\|P g\|_{X}}{\|g\|_{X}} \\
& >\frac{\|P\|_{X \rightarrow X}-\varepsilon\left(1+\|P\|_{X \rightarrow X}\right)}{1+\varepsilon}
\end{aligned}
$$

Hence

$$
c_{X}>\frac{\|P\|_{X \rightarrow X}-\varepsilon\left(1+\|P\|_{X \rightarrow X}\right)}{1+\varepsilon} \quad \text { for all } \varepsilon>0
$$

Passing to the limit as $\varepsilon \rightarrow 0$, we arrive at (6.1).

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## Data availability

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[^0]:    * Corresponding author.

    E-mail addresses: oyk@fct.unl.pt (O. Karlovych), eugene.shargorodsky@kcl.ac.uk (E. Shargorodsky). URLs: https://docentes.fct.unl.pt/oyk (O. Karlovych),
    https://www.kcl.ac.uk/people/eugene-shargorodsky (E. Shargorodsky).

