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# Isometric representation of integral operators with positive-semidefinite kernels 

Bertrand GaUTHIER*


#### Abstract

We describe a natural coisometry from the Hilbert space of all HilbertSchmidt operators on a separable reproducing kernel Hilbert space (RKHS) $\mathcal{H}$, and onto the RKHS $\mathcal{G}$ associated with the squared-modulus of the reproducing kernel of $\mathcal{H}$. We discuss the properties of this coisometry, and show that trace-class integral operators defined by general measures and the reproducing kernel of $\mathcal{H}$ always belong to its initial space. The images of such operators are the Riesz representations of integral functionals on $\mathcal{G}$, drawing a direct connection between the approximation of trace-class integral operators with PSD kernels and the approximation of integral functionals on RKHSs associated with squared-modulus kernels.


Keywords: reproducing kernel Hilbert spaces, integral operators, low-rank approximation, kernel embedding of measures, kernel quadrature.
Mathematics Subject Classification: 46E22, 47B32, 47G10, 65F55.

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## 1 Introduction

Integral operators with positive-semidefinite (PSD) kernels play a central role in the theory of reproducing kernel Hilbert spaces (RKHSs) and their applications; see for instance $[4,5,19,26,20]$. As an important instance, this class of operators encompasses the PSD matrices.

Among other interpretations and under suitable conditions, an integral operator defined by a PSD kernel $K$ and a measure $\mu$ can be regarded as a Hilbert-Schmidt (HS) operator $L_{\mu}$ on the RKHS $\mathcal{H}$ associated with $K$; see e.g. [21, 22]. Let $\mathcal{G}$ be the RKHS for which the squared-modulus kernel $|K|^{2}$ is reproducing. In this note, we show that when the integral of the diagonal of $K$ with respect to the variation of $\mu$ is finite, the HS operator $L_{\mu}$ on $\mathcal{H}$ can be isometrically represented as the Riesz representation $g_{\mu} \in \mathcal{G}$ of the integral functional on $\mathcal{G}$ defined by the measure $\bar{\mu}$, the conjugate of $\mu$. The operator $L_{\mu}$ is in this case trace-class, and $g_{\mu}$ is the kernel embedding, or potential, of the measure $\mu$ in the RKHS $\mathcal{G}$. Consequently, in the Hilbert space $\mathrm{HS}(\mathcal{H})$ of all HS operators on $\mathcal{H}$, the characterisation of an approximate measure $\nu$ (defining together with $K$ an integral operator $L_{\imath}$ on $\mathcal{H}$ ) for the approximation of an initial operator $L_{\mu}$ reduces to a quadrature problem for the approximation of integral functionals on the RKHS $\mathcal{G}$ associated with $|K|^{2}$. More precisely, denoting by $B_{\mathcal{G}}$ the closed unit ball of $\mathcal{G}$, we have

$$
\left\|L_{\mu}-L_{\nu}\right\|_{\mathrm{HS}(\mathcal{H})}=\left\|g_{\mu}-g_{\nu}\right\|_{\mathcal{G}}=\sup _{g \in B_{G}}\left|\int g(t) \mathrm{d} \bar{\mu}(t)-\int g(t) \mathrm{d} \bar{\nu}(t)\right|,
$$

so that the norm $\left\|L_{\mu}-L_{\nu}\right\|_{\mathrm{HS}(\mathcal{H})}$ corresponds to a generalised integral probability metric (see for instance [16, 24, 2, 15, 27]) for the approximation of integral functionals on $\mathcal{G}$.

We give an overall description of the framework surrounding such an isometric representation, and illustrate that it follows from the definition of a natural coisometry $\Gamma$ from $\operatorname{HS}(\mathcal{H})$ onto $\mathcal{G}$; in particular, $\Gamma$ maps self-adjoint operators to real-valued functions, and PSD operators to nonnegative functions (Section 2). Under general measurability conditions on $K$ and assuming that the diagonal of $K$ is integrable with respect to $|\mu|$, we show that the integral operator $L_{\mu}$ always belongs to the initial space of $\Gamma$, and that $\Gamma\left[L_{\mu}\right]=g_{\mu}$. We next describe the equivalence between the quadrature approximation of trace-class integral operators with PSD kernels and kernel quadrature in RKHSs with squared-modulus kernels (Section 3).

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For an approximate measure $\nu$, we denote by $P_{v}$ the orthogonal projection from $\mathcal{H}$ onto $\mathcal{H}_{\nu}$, the closure in $\mathcal{H}$ of the range of $L_{|v|}$. We then discuss the extent to which the kernel-quadrature setting of Section 3 can be used as a surrogate for the characterisation of projection-based approximations of $L_{\mu}$ of the form $P_{\nu} L_{\mu}, L_{\mu} P_{\nu}$ or $P_{\nu} L_{\mu} P_{\nu}$. Further, when the initial measure $\mu$ is nonnegative, the integral operator defined by the kernel $K$ and the measure $\mu$ can more generally be regarded as a HS operator from, and to, $L^{2}(\mu)$ or $\mathcal{H}$ (four possibilities). Noticing that these integral operators are composites of the natural embedding $l_{\mu}$ of $\mathcal{H}$ in $L^{2}(\mu)$, we study the properties of the low-rank approximations induced in this framework by partial embeddings of the form $l_{\mu} P_{v}$. We in parallel illustrate the interest of the kernel-quadrature setting in terms of differentiability and numerical complexity, and describe the connections between of the considered framework and the low-rank approximation of PSD matrices through column sampling (Section 4).

## 2 Framework and notations

By default, all the Hilbert spaces considered in this note are complex; they are otherwise explicitly referred to as real Hilbert spaces; we use a similar convention for vector spaces. Inner products of complex Hilbert spaces are assumed to be linear with respect to their right argument. For $z \in \mathbb{C}$, we denote by $\bar{z},|z|$ and $\Re(z)$ the conjugate, modulus and real part of $z$, respectively, and $\mathrm{i} \in \mathbb{C}$ is the imaginary unit. By analogy, for a complex-valued function $f$ on a general set $S$, we denote by $\bar{f}$ and $|f|$ the functions defined as $\bar{f}(s)=\overline{f(s)}$ and $|f|(s)=|f(s)|, s \in S$; we also use the notation $|f|^{2}$ to refer to the function $s \mapsto|f(s)|^{2}$.

For two Hilbert spaces $H$ and $F$, we denote by $A^{*}$ the adjoint of a bounded linear operator $A: H \rightarrow F$. The map $A$ is an isometry if $A^{*} A=\mathrm{id}_{H}$, the identity operator on $H$, and $A$ is coisometry if $A^{*}$ is an isometry (i.e. if $A A^{*}=\operatorname{id}_{F}$ ). A coisometry $A$ is a surjective partial isometry (i.e. $A A^{*} A=A$ ), and $A^{*} A$ is then the orthogonal projection from $H$ onto the initial space $\mathcal{I}(A)$ of $A$, with $\mathcal{I}(A)$ the orthogonal complement in $H$ of the nullspace of $A$. We denote by $\operatorname{null}(A)$ the nullspace of $A$, and by range $(A)$ its range. Also, for a subset $C$ of $H$, we denote by $C^{\perp_{H}}$ the orthogonal complement of $C$ in $H$, and by $\bar{C}^{H}$ the closure of $C$ in $H$.

### 2.1 Hilbert-Schmidt operators and RKHSs

Hereafter, we introduce the overall framework leading to the definition of a natural coisometry $\Gamma$ from the Hilbert space $\operatorname{HS}(\mathcal{H})$ of all HS operators on an underlying RKHS $\mathcal{H}$, and onto the RKHS $\mathcal{G}$ associated with the squared-modulus of the reproducing kernel of $\mathcal{H}$; the terminology natural is used to emphasise that the considered construction does not depend on the choice of any specific basis.

Underlying RKHS. Let $\mathcal{H}$ be a separable RKHS of complex-valued functions on a general set $\mathscr{X}$, with reproducing kernel $K: \mathscr{X} \times \mathscr{X} \rightarrow \mathbb{C}$; see for instance
[1, 17]. For $t \in \mathscr{X}$, let $k_{t} \in \mathcal{H}$ be defined as $k_{t}(x)=K(x, t)$, with $x \in \mathscr{X}$. For all $h \in \mathcal{H}$, we have $\left\langle k_{t} \mid h\right\rangle_{\mathcal{H}}=h(t)$, where $\langle\cdot \mid \cdot\rangle_{\mathcal{H}}$ stands for the inner product of $\mathcal{H}$ (this equality is often referred to as the reproducing property); we denote by $\|\cdot\|_{\mathcal{H}}$ the norm of $\mathcal{H}$, and we use similar notations for the inner products and norms of all the Hilbert spaces encountered in this note.

Hilbert-Schmidt space. Let $\mathrm{HS}(\mathcal{H})$ be the Hilbert space of all HS operators on $\mathcal{H}$; see e.g. [7, 11, 3]. For $T \in \operatorname{HS}(\mathcal{H})$, we denote by $T[h] \in \mathcal{H}$ the image of $h \in \mathcal{H}$ through $T$, and by $T[h](x)$ the value of the function $T[h]$ at $x \in \mathscr{X}$; we use a similar convention for all function-valued operators. For $a$ and $b \in \mathcal{H}$, let $T_{a, b} \in \operatorname{HS}(\mathcal{H})$ be the rank-one operator given by

$$
T_{a, b}[h]=a\langle b \mid h\rangle_{\mathcal{H}}, h \in \mathcal{H} ;
$$

we also set $S_{b}=T_{b, b}$.
Remark 2.1. An operator $T \in \operatorname{HS}(\mathcal{H})$ admits a singular value decomposition of the form $T=\sum_{i \in \mathbb{\square}} \sigma_{i} T_{u_{i}, v_{i}}, \square \subseteq \mathbb{N}$, where $\left\{\sigma_{i}\right\}_{i \in \mathbb{1}} \in \ell^{2}(\mathbb{\square})$ is the set of all strictly positive singular values of $T$, and where $\left\{u_{i}\right\}_{i \in \cap}$ and $\left\{v_{i}\right\}_{i \in \cap}$ are two orthonormal systems (ONSs) in $\mathcal{H}$; the series converges in $\operatorname{HS}(\mathcal{H})$.

Remark 2.2. Let $\mathcal{H}^{\prime}$ be the continuous dual of $\mathcal{H}$. For $h \in \mathcal{H}$, let $\xi_{h} \in \mathcal{H}^{\prime}$ be the bounded linear functional such that $\xi_{h}(f)=\langle h \mid f\rangle_{\mathcal{H}}, f \in \mathcal{H}$. Endowed with the inner product $\left\langle\xi_{f} \mid \xi_{h}\right\rangle_{\mathcal{H}^{\prime}}=\langle h \mid f\rangle_{\mathcal{H}}$, the vector space $\mathcal{H}^{\prime}$ is a Hilbert space, and the Riesz map $h \mapsto \xi_{h}$ is a bijective conjugate-linear isometry form $\mathcal{H}$ to $\mathcal{H}^{\prime}$ (we may notice that $\xi_{\alpha h}=\bar{\alpha} \xi_{h}, \alpha \in \mathbb{C}$ ). The linear map densely defined as $T_{a, b} \mapsto a \otimes \xi_{b}$ (see Remark 2.1) is then a bijective isometry from the Hilbert space $\mathrm{HS}(\mathcal{H})$ to the tensor Hilbert space $\mathcal{H} \otimes \mathcal{H}^{\prime}$.

Conjugate RKHS. Let $\overline{\mathcal{H}}$ be the RKHS of complex-valued functions on $\mathscr{X}$ associated with the conjugate kernel $\bar{K}$, with $\bar{K}(x, t)=\overline{K(x, t)}, x$ and $t \in \mathscr{X}$. For all $h \in \mathcal{H}$, we have $\bar{h} \in \overline{\mathcal{H}}$ (i.e. the function $\bar{h}: x \mapsto \overline{h(x)}$ is a vector of $\overline{\mathcal{H}}$ ), and the map $h \mapsto \bar{h}$ is a bijective conjugate-linear isometry from $\mathcal{H}$ to $\overline{\mathcal{H}}$. We may notice that $\overline{k_{t}}(x)=\overline{K(x, t)}=\overline{k_{t}(x)}$, and that

$$
\left\langle\overline{k_{x}} \mid \overline{k_{t}}\right\rangle_{\overline{\mathcal{H}}}=\overline{K(x, t)}=K(t, x)=\left\langle k_{t} \mid k_{x}\right\rangle_{\mathcal{H}} .
$$

We denote by $\Psi$ the bijective linear isometry from $\operatorname{HS}(\mathcal{H})$ to the tensor Hilbert space $\mathcal{H} \otimes \overline{\mathcal{H}}$, densely defined as $\Psi\left(T_{a, b}\right)=a \otimes \bar{b}, a$ and $b \in \mathcal{H}$.

Remark 2.3. Following Remark 2.2, the linear map $\xi_{h} \mapsto \bar{h}$ is a bijective isometry form $\mathcal{H}^{\prime}$ to $\overline{\mathcal{H}}$. Further, the linear map densely defined as $a \otimes \xi_{b} \mapsto a \otimes \bar{b}$ is a bijective isometry form $\mathcal{H} \otimes \mathcal{H}^{\prime}$ to $\mathcal{H} \otimes \overline{\mathcal{H}}$; the composition of this isometry with the bijective isometry from $\operatorname{HS}(\mathcal{H})$ to $\mathcal{H} \otimes \mathcal{H}^{\prime}$ discussed in Remark 2.2 yields the isometry $\Psi: \operatorname{HS}(\mathcal{H}) \rightarrow \mathcal{H} \otimes \overline{\mathcal{H}}$; see the diagram (5).

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Squared-kernel RKHS. The kernels $K$ and $\bar{K}$ being PSD, by the Schur-product theorem, so is the squared-modulus kernel $|K|^{2}=K \bar{K}$, with

$$
|K|^{2}(x, t)=K(x, t) \overline{K(x, t)}=|K(x, t)|^{2}=\left|k_{t}\right|^{2}(x), x \text { and } t \in \mathscr{X} .
$$

Let $\mathcal{G}$ be the RKHS of complex-valued functions on $\mathscr{X}$ for which $|K|^{2}$ is reproducing ( $\mathcal{G}$ is the product of the two RKHSs $\mathcal{H}$ and $\overline{\mathcal{H}}$; see e.g. [1, 17]).

Following [17, Chapter 5], we denote by $C_{\Delta}: \mathcal{H} \otimes \overline{\mathcal{H}} \rightarrow \mathcal{G}$ the coisometry densely defined as $C_{\Delta}(a \otimes \bar{b})=a \bar{b}, a$ and $b \in \mathcal{H}$, where $a \bar{b} \in \mathcal{G}$ is the complexvalued function on $\mathscr{X}$ given by

$$
(a \bar{b})(x)=a(x) \overline{b(x)}=\left\langle k_{x} \otimes \overline{k_{x}} \mid a \otimes \bar{b}\right\rangle_{\mathcal{H} \otimes \overline{\mathcal{H}}}, x \in \mathscr{X} .
$$

For $\Upsilon \in \mathcal{H} \otimes \mathcal{H}^{\prime}$, we more generally have

$$
\begin{equation*}
C_{\Delta}[\Upsilon](x)=\left.\langle | k_{x}\right|^{2}\left|C_{\Delta}[\Upsilon]\right\rangle_{\mathcal{G}}=\left\langle k_{x} \otimes \overline{k_{x}} \mid \Upsilon\right\rangle_{\mathcal{H} \otimes \overline{\mathcal{H}}} . \tag{1}
\end{equation*}
$$

The initial space of $C_{\Delta}$ is $\mathcal{I}\left(C_{\Delta}\right)=\overline{\operatorname{span}_{\mathbb{C}}\left\{k_{x} \otimes \overline{k_{x}} \mid x \in \mathscr{X}\right\}}{ }^{\mathcal{H} \otimes \mathcal{H}}$, the closure in $\mathcal{H} \otimes \overline{\mathcal{H}}$ of the linear space spanned by the simple tensors $k_{x} \otimes \overline{k_{x}}, x \in \mathscr{X}$.

Remark 2.4. From (1), for all $x \in \mathscr{X}$, we have $C_{\Delta}^{*}\left[\left|k_{x}\right|^{2}\right]=k_{x} \otimes \overline{k_{x}}$. The linear space $\operatorname{span}_{\mathbb{C}}\left\{\left|k_{x}\right|^{2} \mid x \in \mathscr{X}\right\}$ being dense in $\mathcal{G}$ (see for instance [17, Chapter 2]), we have $C_{\Delta} C_{\Delta}^{*}=\mathrm{id}_{C}$, and $C_{\Delta}^{*}$ is thus an isometry.

Natural coisometry from $\operatorname{HS}(\mathcal{H})$ onto $\mathcal{G}$. We are now in the position to introduce the coisometry $\Gamma=C_{\Delta} \Psi: \operatorname{HS}(\mathcal{H}) \rightarrow \mathcal{G}$; since this map plays a central role in our study, we gather some of its properties in the following Lemma 2.1.

Lemma 2.1. The linear map $\Gamma=C_{\Delta} \Psi$ is a coisometry from $\operatorname{HS}(\mathcal{H})$ onto $\mathcal{G}$, and its initial space verifies

$$
\begin{equation*}
\mathcal{I}(\Gamma)=\overline{\operatorname{span}}_{\mathbb{C}}\left\{S_{k_{x}} \mid x \in \mathscr{X}\right\} \quad \mathrm{HS}(\mathcal{H}) \tag{2}
\end{equation*}
$$

Further, for all $T \in \operatorname{HS}(\mathcal{H})$, we have

$$
\begin{equation*}
\Gamma[T](x)=\left\langle S_{k_{x}} \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})}=\left\langle k_{x} \mid T\left[k_{x}\right]\right\rangle_{\mathcal{H}}=T\left[k_{x}\right](x), x \in \mathscr{X} . \tag{3}
\end{equation*}
$$

Proof. The linear isometry $\Psi$ being bijective, we have $\Psi \Psi^{*}=\mathrm{id}_{\mathcal{H} \otimes \overline{\mathcal{H}}}$, and thus

$$
\Gamma \Gamma^{*}=C_{\Delta} \Psi \Psi^{*} C_{\Delta}^{*}=C_{\Delta} C_{\Delta}^{*}=\mathrm{id}_{\mathcal{G}} .
$$

By definition of $C_{\Delta}$ and $\Psi$, we have $\Gamma\left[T_{a, b}\right]=C_{\Delta}[a \otimes \bar{b}]=a \bar{b}, a$ and $b \in \mathcal{H}$. Noticing that $\Gamma^{*}\left[\left|k_{x}\right|^{2}\right]=\Psi^{*}\left[k_{x} \otimes \overline{k_{x}}\right]=S_{k_{x}}, x \in \mathscr{X}$, (2) follows by density of $\operatorname{span}_{\mathbb{C}}\left\{\left|k_{x}\right|^{2} \mid x \in \mathscr{X}\right\}$ in $\mathcal{G}$ (see Remark 2.4). From the reproducing property in $\mathcal{G}$,
we have $\Gamma[T](x)=\left.\langle | k_{x}\right|^{2}|\Gamma[T]\rangle_{\mathcal{C}}, T \in \operatorname{HS}(\mathcal{H})$; as $\Gamma^{*}\left[\left|k_{x}\right|^{2}\right]=S_{k_{x}}$, we obtain $\Gamma[T](x)=\left\langle S_{k_{x}} \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})}$. We next notice that

$$
\begin{equation*}
\left\langle T_{a, b} \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})}=\langle a \mid T[b]\rangle_{\mathcal{H}}, a \text { and } b \in \mathcal{H} \tag{4}
\end{equation*}
$$

indeed, as $T_{a, 0}=0$, equality (4) trivially holds for $b=0$, and for $b \neq 0$, we have $\left\langle T_{a, b} \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})}=\left\langle T_{a, b}[b] \mid T[b]\right\rangle_{\mathcal{H}} /\|b\|_{\mathcal{H}}^{2}$, with $T_{a, b}[b]=a\|b\|_{\mathcal{H}}^{2}$. Taking $a=b=k_{x}$, (4) reads $\left\langle S_{k_{x}} \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})}=\left\langle k_{x} \mid T\left[k_{x}\right]\right\rangle_{\mathcal{H}}=T\left[k_{x}\right](x)$, concluding the proof.

The following diagram summarises the construction of $\Gamma$ (the $\cong$ symbol refers to the two the natural bijective isometries discussed in Remarks 2.2 and 2.3).


Through $\Gamma$, the HS operators on $\mathcal{H}$ belonging to $\mathcal{I}(\Gamma)$ are isometrically represented as functions in the RKHS $\mathcal{G}$ associated with the squared-modulus kernel $|K|^{2}$.

### 2.2 Self-adjoint operators and real-valuedness

For $T \in \operatorname{HS}(\mathcal{H})$, we have $T^{*} \in \operatorname{HS}(\mathcal{H})$, and so, from (3), $\Gamma\left[T^{*}\right]=\overline{\Gamma[T]}$. Consequently, if $T$ is self-adjoint, then the function $\Gamma[T]$ is real-valued. This observation is intrinsically related to the structures of $\operatorname{HS}(\mathcal{H})$ and $\mathcal{G}$, as illustrated below.

Remark 2.5. Form the surjectivity of $\Gamma$, for all $g \in \mathcal{G}$, there exists $T \in \operatorname{HS}(\mathcal{H})$ such that $\Gamma[T]=g$, and the function $\bar{g}=\Gamma\left[T^{*}\right]$ is thus also a vector of $\mathcal{G}$.

Squared-kernel real RKHS. The squared kernel $|K|^{2}$ being real-valued, it is the reproducing kernel of a real RKHS $\mathcal{G}_{\mathbb{R}}$ of real-valued functions on $\mathscr{X}$. The real RKHS $\mathcal{G}_{\mathbb{R}}$ can be regarded as a closed real-linear subspace of $\mathcal{G}$, and the complex RKHS $\mathcal{G}$ is a complexification of $\mathcal{G}_{\mathbb{R}}$ (see e.g. [17, Chapter 5]; $\mathcal{G}_{\mathbb{R}}$ is referred to as a standard subspace of $\mathcal{G}$ ). More precisely, for all $g \in \mathcal{G}$, there exist $g_{\mathrm{r}}$ and $g_{\mathrm{i}} \in \mathcal{G}_{\mathbb{R}} \subset \mathcal{G}$ such that $g=g_{\mathrm{r}}+\mathrm{i} g_{\mathrm{i}}$, with $g_{\mathrm{r}}=\frac{1}{2}(\bar{g}+g)$ and $g_{\mathrm{i}}=\frac{1}{2} \mathrm{i}(\bar{g}-g)$, and the decomposition $\mathcal{G}=\mathcal{G}_{\mathbb{R}}+\mathrm{i} \mathcal{G}_{\mathbb{R}}$ holds. The functions $g_{\mathrm{r}}$ and $g_{\mathrm{i}}$ are the real and imaginary parts of $g$, and we have $\|g\|_{\mathcal{G}}^{2}=\left\|g_{r}\right\|_{\mathcal{G}}^{2}+\left\|g_{\mathrm{i}}\right\|_{\mathcal{G}}^{2}$.

Self-adjoint HS operators. We denote by $\operatorname{HS}_{*}(\mathcal{H}) \subset \operatorname{HS}(\mathcal{H})$ the closed real-linear subspace of all self-adjoint HS operators on $\mathcal{H}$. For every $T \in \operatorname{HS}(\mathcal{H})$, there exist $T_{r}$ and $T_{\mathrm{i}} \in \mathrm{HS}_{*}(\mathcal{H})$ such that $T=T_{\mathrm{r}}+\mathrm{i} T_{\mathrm{i}}$, with $T_{\mathrm{r}}=\frac{1}{2}\left(T^{*}+T\right)$ and $T_{\mathrm{i}}=\frac{1}{2} \mathrm{i}\left(T^{*}-T\right)$. Endowed with the inner product of $\operatorname{HS}(\mathcal{H})$, the real vector space $\operatorname{HS}_{*}(\mathcal{H})$ is a real Hilbert space, and the decomposition $\operatorname{HS}(\mathcal{H})=\operatorname{HS}_{*}(\mathcal{H})+\mathrm{iHS}_{*}(\mathcal{H})$ holds. By linearity of $\Gamma$, we may notice that if $\Gamma[T]=g$, then $\Gamma\left[T_{\mathrm{r}}\right]=g_{\mathrm{r}}$ and $\Gamma\left[T_{\mathrm{i}}\right]=g_{\mathrm{i}}$.

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Remark 2.6. For all $T \in \operatorname{HS}(\mathcal{H})$, we have $\|T\|_{\mathrm{HS}(\mathcal{H})} \geqslant\|\Gamma[T]\|_{\mathcal{G}}$, with equality if and only if $T \in \mathcal{I}(\Gamma)$; as $\|T\|_{\mathrm{HS}(\mathcal{H})}=\left\|T^{*}\right\|_{\mathrm{HS}(\mathcal{H})}$ and $\|g\|_{\mathcal{G}}=\|\bar{g}\|_{G}$, it follows that if $T \in \mathcal{I}(\Gamma)$, then $T^{*} \in \mathcal{I}(\Gamma)$, and also $T_{\mathrm{r}}$ and $T_{\mathrm{i}} \in \mathcal{I}(\Gamma)$.

Remark 2.7. The diagram (5) is also well-defined when all the involved Hilbert spaces are real (the kernel $K$ in this case verifies $K=\bar{K}$, that is, $K$ is symmetric). We then have $\mathcal{I}(\Gamma)=\overline{\operatorname{span}}_{\mathbb{R}}\left\{S_{k_{x}} \mid x \in \mathscr{X}\right\} \quad$ 'HS $(\mathcal{H})$, and $\mathcal{I}(\Gamma)$ thus solely consists of selfadjoint operators. In the real case, as $\Gamma\left[T^{*}\right]=\Gamma[T], T \in \operatorname{HS}(\mathcal{H})$, we also obtain that if $T^{*}=-T$, then $\Gamma[T]=0$; by comparison, in the complex case, if $T^{*}=-T$, then the function $\Gamma[T]$ is pure-imaginary.

### 2.3 Positive-semidefinite operators and nonnegativity

From (3), if $T \in \operatorname{HS}(\mathcal{H})$ is PSD, then $\Gamma[T](x)=\left\langle k_{x} \mid T\left[k_{x}\right]\right\rangle_{\mathcal{H}} \geqslant 0$ for all $x \in \mathscr{X}$, so that the function $\Gamma[T]$ is nonnegative. This observation again directly relates to the structures of $\mathrm{HS}(\mathcal{H})$ and $\mathcal{G}$.

Positive-semidefinite $H S$ operators. We denote by $\operatorname{HS}_{*}^{+}(\mathcal{H}) \subset \mathrm{HS}_{*}(\mathcal{H})$ the closed real convex cone of all PSD HS operators on $\mathcal{H}$. For all $T \in \mathrm{HS}_{*}(\mathcal{H})$, there exist $T^{+}$ and $T^{-} \in \operatorname{HS}_{*}^{+}(\mathcal{H})$ such that $T=T^{+}-T^{-}$, so that the cone $\operatorname{HS}_{*}^{+}(\mathcal{H})$ is generating in $\mathrm{HS}_{*}(\mathcal{H})$.

Nonnegative functions in $\mathcal{G}$. We denote by $\mathcal{C}_{\mathbb{R}}^{+} \subset \mathcal{G}_{\mathbb{R}}$ the closed real convex cone of all nonnegative functions in $\mathcal{G}_{\mathbb{R}}$. For all $T \in \mathrm{HS}_{*}^{+}(\mathcal{H})$, we have $\Gamma[T] \in \mathcal{G}_{\mathbb{R}}^{+}$, and since $\Gamma$ is surjective, the cone $\mathcal{C}_{\mathbb{R}}^{+}$is generating (that is, for all $g \in \mathcal{G}_{\mathbb{R}}$, there exist $g^{+}$ and $g^{-} \in \mathcal{G}_{\mathbb{R}}^{+}$such that $g=g^{+}-g^{-}$).

Remark 2.8. Let $T=\sum_{j \in \|} \lambda_{j} S_{\varphi_{j}}$ be a spectral expansion of $T \in \operatorname{HS}_{*}^{+}(\mathcal{H})$, with $\left\{\lambda_{j}\right\}_{j \in 1}, \mathbb{\square} \subseteq \mathbb{N}$, the set of all strictly-positive eigenvalues of $T$ (repeated with multiplicity), and where $\left\{\varphi_{j}\right\}_{j \in \mathbb{1}} \subset \mathcal{H}$ is a set of associated eigenvectors, orthonormal in $\mathcal{H}$. We then have $\Gamma[T]=\sum_{j \in ⿺} \lambda_{j} \Gamma\left[S_{\varphi_{j}}\right]=\sum_{j \in 0} \lambda_{j}\left|\varphi_{j}\right|^{2}$, so that if $T \in \operatorname{HS}(\mathcal{H})$ is PSD, then $\Gamma[T]=0$ if and only if $T=0$.

## 3 Trace-class integral operators with PSD kernels

From Lemma 2.1, if $T \in \operatorname{HS}(\mathcal{H})$ is of the form $T=\sum_{j=1}^{n} \omega_{j} S_{k_{s_{j}}}$, with $n \in \mathbb{N}$, $s_{j} \in \mathscr{X}$ and $\omega_{j} \in \mathbb{C}$, then $T \in \mathcal{I}(\Gamma)$. For $h \in \mathcal{H}$, we in this case have

$$
T[h](x)=\sum_{j=1}^{n} \omega_{j} k_{s_{j}}(x)\left\langle k_{s_{j}} \mid h\right\rangle_{\mathcal{H}}=\sum_{j=1}^{n} \omega_{j} K\left(x, s_{j}\right) h\left(s_{j}\right), x \in \mathscr{X},
$$

so that $T$ can be regarded as an integral operator on $\mathcal{H}$ defined by the kernel $K$ and the finitely-supported measure $\sum_{j=1}^{n} \omega_{j} \delta_{s_{j}}$, with $\delta_{x}$ the Dirac measure at $x \in \mathscr{X}$.

We in addition have $\Gamma[T]=\sum_{j=1}^{n} \omega_{j}\left|k_{s_{j}}\right|^{2}$, so that $\langle\Gamma[T] \mid g\rangle_{\mathcal{G}}=\sum_{j=1}^{n} \bar{\omega}_{j} g\left(s_{j}\right)$, $g \in \mathcal{G}$, and $\Gamma[T]$ is thus the Riesz representation of the integral functional on $\mathcal{G}$ defined by the measure $\sum_{j=1}^{n} \overline{\omega_{j}} \delta_{s_{i}}$. Under adequate measurability conditions, these observations more generally hold for all trace-class integral operators on $\mathcal{H}$ defined by the reproducing kernel $K$ of $\mathcal{H}$ and general measures on $\mathscr{X}$.

### 3.1 Integral operators and kernel embeddings of measures

Let $\mathcal{A}$ be a $\sigma$-algebra of subsets of $\mathscr{X}$. We consider the Borel $\sigma$-algebra of $\mathbb{C}$, and make the following assumptions on $K$ and the measurable space $(\mathscr{X}, \mathcal{A})$ :
(A.1) for all $t \in \mathscr{X}$, the function $k_{t} \in \mathcal{H}$ is measurable on $(\mathscr{X}, \mathcal{A})$;
(A.2) the diagonal of $K$ is measurable on $(\mathscr{X}, \mathcal{A})$.

Remark 3.1. The RKHSs $\mathcal{H}$ and $\mathcal{G}$ being separable, (A.1) ensures that all the functions in $\mathcal{H}$ and $\mathcal{C}$ are measurable; see for instance [25, Lemma 4.24]. The three maps $t \mapsto k_{t}, t \mapsto\left|k_{t}\right|^{2}$ and $t \mapsto S_{k_{t}}, t \in \mathscr{X}$, are then weakly-measurable, and since the Hilbert spaces $\mathcal{H}, \mathcal{G}$ and $\operatorname{HS}(\mathcal{H})$ are separable, by the Pettis measurability theorem, these three maps are also strongly-measurable (see e.g. [8, 28]).

We denote by $\mathcal{M}_{+}, \mathcal{M}$, and $\mathcal{M}_{\mathbb{C}}$ the set of all nonnegative, signed and complex measures ${ }^{1}$ on $(\mathscr{X}, \mathcal{A})$, and we set $\mathcal{M}_{\mathbb{F}}=\mathcal{M} \cup \mathcal{M}_{\mathbb{C}}$ (we have $\mathcal{M}_{+} \subset \mathcal{M}$ ). Noticing that $K(t, t) \geqslant 0, t \in \mathscr{X}$, from (A.2), we define

$$
\tau_{\mu}=\int_{\mathscr{X}} K(t, t) \mathrm{d}|\mu|(t) \in \mathbb{R}_{\geqslant 0} \cup\{+\infty\}, \mu \in \mathcal{M}_{\mathbb{F}} .
$$

We next introduce the sets $\mathcal{T}_{+}(K), \mathcal{T}(K)$ and $\mathcal{T}_{\mathbb{C}}(K)$ of all measures $\mu$ in $\mathcal{M}_{+}, \mathcal{M}$, and $\mathcal{M}_{\mathbb{C}}$ such that $\tau_{\mu}$ is finite, respectively; the inclusion $\mathcal{T}_{+}(K) \subset \mathcal{T}(K)$ holds, and we set $\mathcal{T}_{\mathbb{F}}(K)=\mathcal{T}(K) \cup \mathcal{T}_{\mathbb{C}}(K)$. We may notice that $\mathcal{T}_{\mathbb{C}}(K)$ is a vector space, and that $\mathcal{T}_{+}(K)$ is a real convex cone. We recall that

$$
K(t, t)=\left\|k_{t}\right\|_{\mathcal{H}}^{2}=\left\|S_{k_{t}}\right\|_{\mathrm{HS}(\mathcal{H})}=\left\|\left|k_{t}\right|^{2}\right\|_{\mathcal{G}}, t \in \mathscr{X} .
$$

Integral operators on $\mathcal{H}$ with kernel $K$. By assumption, for $\mu \in \mathcal{T}_{\mathbb{F}}(K)$, the integral $\int_{\mathscr{X}}\left\|S_{k_{t}}\right\|_{\mathrm{HS}(\mathcal{H})} \mathrm{d}|\mu|(t)=\tau_{\mu}$ is finite, and the map $t \mapsto S_{k_{t}}$ is thus Bochner-integrable with respect to $\mu$ (Bochner integrability criterion; see e.g. [8, 28]). We then set

$$
L_{\mu}=\int_{\mathscr{X}} S_{k_{t}} \mathrm{~d} \mu(t) \in \operatorname{HS}(\mathcal{H}) .
$$

From (4), for $h \in \mathcal{H}$ and $x \in \mathscr{X}$, we have

$$
L_{\mu}[h](x)=\left\langle T_{k_{x}, h} \mid L_{\mu}\right\rangle_{\mathrm{HS}(\mathcal{H})}=\int_{\mathscr{X}}\left\langle T_{k_{x}, h} \mid S_{k_{t}}\right\rangle_{\mathrm{HS}(\mathcal{H})} \mathrm{d} \mu(t)=\int_{\mathscr{X}} K(x, t) h(t) \mathrm{d} \mu(t),
$$

and $L_{\mu} \in \operatorname{HS}(\mathcal{H})$ can thus be regarded as an integral operator on $\mathcal{H}$ defined by the kernel $K$ and the measure $\mu$ (see also Remark 3.2).

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Lemma 3.1. For all $\mu \in \mathcal{T}_{\mathfrak{F}}(K)$, we have

$$
\begin{equation*}
\left\langle L_{\mu}[h] \mid f\right\rangle_{\mathcal{H}}=\int_{\mathscr{X}} \overline{h(t)} f(t) \mathrm{d} \bar{\mu}(t), h \text { and } f \in \mathcal{H} . \tag{6}
\end{equation*}
$$

Proof. From the Cauchy-Schwarz (CS) inequality in $\mathcal{H}$, we have

$$
\left|\int_{\mathscr{X}} \overline{h(t)} f(t) \mathrm{d} \bar{\mu}(t)\right| \leqslant \int_{\mathscr{X}}\left|\left\langle k_{t} \mid h\right\rangle_{\mathcal{H}}\left\|\left\langle k_{t} \mid f\right\rangle_{\mathcal{H}}|\mathrm{d}| \mu \mid(t) \leqslant\right\| f\left\|_{\mathcal{H}}\right\| h \|_{\mathcal{H}} \tau_{\mu},\right.
$$

so that the linear functional $\Theta_{h, \mu}: f \mapsto \int_{\mathscr{X}} \overline{h(t)} f(t) \mathrm{d} \bar{\mu}(t), f \in \mathcal{H}$, is bounded. Taking $f=k_{x}, x \in \mathscr{X}$, we obtain

$$
\Theta_{h, \mu}\left(k_{x}\right)=\int_{\mathscr{X}} \overline{h(t)} K(t, x) \mathrm{d} \bar{\mu}(t)=\overline{L_{\mu}[h](x)}=\left\langle L_{\mu}[h] \mid k_{x}\right\rangle_{\mathcal{H}},
$$

so that $L_{\mu}[h]$ is the Riesz representation of $\Theta_{h, \mu}$.
Remark 3.2. For $\mu \in \mathcal{T}_{\mathbb{F}}(K)$, the operator $L_{\mu} \in \operatorname{HS}(\mathcal{H})$ is the Riesz representation of the bounded linear functional $Z_{\mu}: \operatorname{HS}(\mathcal{H}) \rightarrow \mathbb{C}$ given by

$$
Z_{\mu}(T)=\int_{\mathscr{X}}\left\langle S_{k_{t}} \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})} \mathrm{d} \bar{\mu}(t), T \in \operatorname{HS}(\mathcal{H}),
$$

that is $Z_{\mu}(T)=\left\langle L_{\mu} \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})}$. Using the CS inequality in $\operatorname{HS}(\mathcal{H})$, we in particular have $\left|Z_{\mu}(T)\right| \leqslant \int_{\mathscr{X}}\left|\left\langle S_{k_{t}} \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})}\right| \mathrm{d}|\mu|(t) \leqslant\|T\|_{\mathrm{HS}(\mathcal{H})} \tau_{\mu}$.

We may also notice that the map $t \mapsto k_{t}$ is Bochner-integrable with respect to the measure $\mathfrak{m}_{h, \mu}$ on $(\mathscr{X}, \mathcal{A})$ defined as $\mathfrak{m}_{h, \mu}(A)=\int_{A} h(t) \mathrm{d} \mu(t), A \in \mathcal{A}$; we then have $L_{\mu}[h]=\int_{\mathscr{X}} k_{t} h(t) \mathrm{d} \mu(t)$, with $\int_{\mathscr{X}}\left\|k_{t}\right\|_{\mathcal{H}}|h(t)| \mathrm{d}|\mu|(t) \leqslant\|h\|_{\mathcal{H}} \tau_{\mu}$.

Lemma 3.2. For all $\mu \in \mathcal{J}_{\mathbb{F}}(K)$, the operator $L_{\mu}$ is trace-class.
Proof. Let $\left\{h_{i}\right\}_{i \in \square}$ be an orthonormal basis (ONB) of $\mathcal{H}$, with $\mathbb{\square} \subseteq \mathbb{N}$. For all $t \in \mathscr{X}$, we have $k_{t}=\sum_{i \in \mathbb{\square}} h_{i} \overline{h_{i}(t)}$, so that $\left\{h_{i}(t)\right\}_{i \in \mathbb{\square}} \in \ell^{2}(\mathbb{\square})$ and $\sum_{i \in \mathbb{I}}\left|h_{i}(t)\right|^{2}=K(t, t)$; see e.g. [17, Chapter 2]. Let $\left\{f_{i}\right\}_{i \in 1}$ be another ONB of $\mathcal{H}$; from (6), and by monotone convergence and the CS inequality in $\ell^{2}(\mathbb{\square})$, we have

$$
\begin{aligned}
\sum_{i \in \mathrm{I}} \mid\left\langle f_{i}\right| & \left.L_{\mu}\left[h_{i}\right]\right\rangle_{\mathcal{H}}\left|\leqslant \sum_{i \in \mathbb{\square}} \int_{\mathscr{X}}\right| f_{i}(t)| | h_{i}(t)|\mathrm{d}| \mu \mid(t) \\
& =\int_{\mathscr{X}} \sum_{i \in \mathbb{V}}\left|f_{i}(t)\right|\left|h_{i}(t)\right| \mathrm{d}|\mu|(t) \leqslant \int_{\mathscr{X}} \sqrt{K(t, t)} \sqrt{K(t, t)} \mathrm{d}|\mu|(t)=\tau_{\mu},
\end{aligned}
$$

so that $\operatorname{trace}\left(\left|L_{\mu}\right|\right) \leqslant \tau_{\mu}$, with $\left|L_{\mu}\right|=\left(L_{\mu}^{*} L_{\mu}\right)^{1 / 2}$ the modulus of $L_{\mu}$.

Kernel embedding of measures in $\mathcal{G}$. As $\int_{\mathscr{X}}\left\|\left|k_{t}\right|^{2}\right\|_{\mathcal{G}} \mathrm{d}|\mu|(t)=\tau_{\mu}$, for $\mu \in \mathcal{T}_{\mathbb{F}}(K)$, the map $t \mapsto\left|k_{t}\right|^{2}$ is Bochner-integrable with respect to $\mu$. We then set

$$
g_{\mu}=\int_{\mathscr{X}}\left|k_{t}\right|^{2} \mathrm{~d} \mu(t) \in \mathcal{C} .
$$

For all $g \in \mathcal{G}$, we have $\left\langle g_{\mu} \mid g\right\rangle_{\mathcal{G}}=\int_{\mathscr{X}} g(t) \mathrm{d} \bar{\mu}(t)$, and $g_{\mu}$ is thus the Riesz representation of the linear functional $I_{\mu}: \mathcal{G} \rightarrow \mathbb{C}$ such that $I_{\mu}(g)=\int_{\mathscr{X}} g(t) \mathrm{d} \bar{\mu}(t)$; we may notice that $\left|I_{\mu}(g)\right| \leqslant\|g\|_{\mathcal{G}} \tau_{\mu}$. The vector $g_{\mu}$ is referred to as the kernel embedding, or potential, of the measure $\mu$ in the RHKS $\mathcal{G}$; see for instance $[24,6,15]$.

Theorem 3.1. For all $\mu \in \mathcal{T}_{\mathbb{F}}(K)$, we have $L_{\mu} \in \mathcal{I}(\Gamma)$ and $\Gamma\left[L_{\mu}\right]=g_{\mu}$.
Proof. From Lemma 2.1 and by definition of $L_{\mu}$ and $g_{\mu}$, for all $T \in \operatorname{HS}(\mathcal{H})$, we have

$$
\left\langle\Gamma^{*}\left[g_{\mu}\right] \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})}=\left\langle g_{\mu} \mid \Gamma[T]\right\rangle_{\mathcal{G}}=\int_{\mathscr{X}}\left\langle S_{k_{t}} \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})} \mathrm{d} \bar{\mu}(t)=\left\langle L_{\mu} \mid T\right\rangle_{\mathrm{HS}(\mathcal{H})},
$$

so that $\Gamma^{*}\left[g_{\mu}\right]=L_{\mu}$.
Remark 3.3. Following Section 2.2, for a signed measure $\mu \in \mathcal{T}(K)$, the function $g_{\mu}$ is real-valued, and the operator $L_{\mu}=\Gamma^{*}\left[g_{\mu}\right]$ is self-adjoint; as expected. Also, following Section 2.3, for a nonnegative measure $\mu \in \mathcal{T}_{+}(K)$, the function $g_{\mu}$ is nonnegative, and the operator $L_{\mu}$ is PSD. We may notice that $L_{\delta_{x}}=S_{k_{x}}, x \in \mathscr{X} . \triangleleft$

### 3.2 Approximation and kernel quadrature

Let $B_{\mathcal{G}}=\left\{g \in \mathcal{G} \mid\|g\|_{\mathcal{G}} \leqslant 1\right\}$ be the closed unit ball of $\mathcal{G}$. For $\mu$ and $v \in \mathcal{T}_{\mathbb{F}}(K)$, from the definition of $g_{\mu}$ and $g_{\nu} \in \mathcal{G}$ and by the CS inequality, we have

$$
\begin{aligned}
\mathfrak{M}_{\mathcal{G}}(\mu, \nu)=\sup _{g \in B_{\mathcal{G}}}\left|\left\langle g_{\mu}-g_{\nu} \mid g\right\rangle_{\mathcal{G}}\right| & =\sup _{g \in B_{\mathcal{G}}}\left|\int_{\mathscr{X}} g(t) \mathrm{d} \bar{\mu}(t)-\int_{\mathscr{X}} g(t) \mathrm{d} \bar{\nu}(t)\right| \\
& =\left\|g_{\mu}-g_{\nu}\right\|_{\mathcal{G}} .
\end{aligned}
$$

When the measures $\mu$ and $v$ are probability measures, the term $\mathfrak{M}_{\mathcal{G}}(\mu, \nu)$ is referred to as an integral probability metric, or maximum mean discrepancy; see for instance [16, 24, 23, 15, 27].

Remark 3.4. For all $g \in \mathcal{G}$, we have $\bar{g} \in \mathcal{G}$ and $\|\bar{g}\|_{\mathcal{G}}=\|g\|_{\mathcal{G}}$ (see Section 2.2), so that $\mathfrak{M}_{\mathcal{G}}(\mu, \nu)=\sup _{g \in B_{G}}\left|\int_{\mathscr{X}} g(t) \mathrm{d} \mu(t)-\int_{\mathscr{X}} g(t) \mathrm{d} v(t)\right|, \mu$ and $v \in \mathcal{T}_{\mathbb{F}}(K)$.

The following Corollary 3.1 illustrates the connection between the quadrature approximation of trace-class integral operators with PSD kernels (quadrature method; see e.g. [13]) and the approximation of integral functionals on RKHSs with squaredmodulus kernels (kernel quadrature; see e.g. [2]).

Corollary 3.1. For all $\mu$ and $\nu \in \mathcal{T}_{\mathbb{F}}(K)$, we have $\left\|L_{\mu}-L_{v}\right\|_{\mathrm{HS}(\mathcal{H})}=\mathfrak{M}_{G}(\mu, \nu)$.

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Proof. Form Theorem 3.1 and by linearity of $\Gamma^{*}$, we have $\Gamma^{*}\left[g_{\mu}-g_{\nu}\right]=L_{\mu}-L_{\nu}$. Since $\Gamma^{*}$ is an isometry, it follows that $\left\|g_{\mu}-g_{\nu}\right\|_{\mathcal{G}}=\left\|L_{\mu}-L_{\nu}\right\|_{\mathrm{HS}(\mathcal{H})}$.

For a fixed initial measure $\mu \in \mathcal{T}_{\mathbb{F}}(K)$, the map $D_{\mu}: \mathcal{J}_{\mathbb{F}}(K) \rightarrow \mathbb{R}_{\geqslant 0}$, with

$$
\begin{equation*}
D_{\mu}(v)=\left\|L_{\mu}-L_{\nu}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=\left\|g_{\mu}-g_{\nu}\right\|_{G}^{2}, v \in \mathcal{T}_{\mathbb{F}}(K), \tag{7}
\end{equation*}
$$

is convex on any convex set $\mathscr{C} \subseteq \mathcal{T}_{\mathbb{F}}(K)$; the map $v \mapsto g_{\nu}$ is indeed linear on $\mathscr{C}$ (that is, for $\xi=v+\alpha \eta \in \mathscr{C}$, with $v, \eta \in \mathscr{C}$ and $\alpha \in \mathbb{C}$, we have $g_{\xi}=g_{v}+\alpha g_{\eta}$ ), and the map $g \mapsto\left\|g_{\mu}-g\right\|_{\mathcal{G}}^{2}$ is convex on $\mathcal{G}$. Also, for $v$ and $\eta \in \mathscr{C}$, the directional derivative of $D_{\mu}$ at $v$ in the direction $\eta-v$ is

$$
\lim _{\rho \rightarrow 0^{+}} \frac{1}{\rho}\left[D_{\mu}(v+\rho(\eta-v))-D_{\mu}(v)\right]=2 \Re\left(\left\langle g_{v}-g_{\mu} \mid g_{\eta}-g_{\nu}\right\rangle_{\mathcal{G}}\right) .
$$

From Theorem 3.1, we may in particular notice that

$$
\begin{align*}
\left\langle g_{v} \mid g_{\mu}\right\rangle_{\mathcal{G}} & =\left\langle L_{v} \mid L_{\mu}\right\rangle_{\mathrm{HS}(\mathcal{H})}=\iint_{\mathscr{X}}|K(x, t)|^{2} \mathrm{~d} \mu(t) \mathrm{d} \bar{v}(x) \\
& =\int_{\mathscr{X}} g_{\mu}(x) \mathrm{d} \bar{v}(x)=\int_{\mathscr{X}} \overline{g_{v}(t)} \mathrm{d} \mu(t) \tag{8}
\end{align*}
$$

and since $|K(x, t)|^{2}=\left|\left\langle k_{x} \mid k_{t}\right\rangle_{\mathcal{H}}\right|^{2} \leqslant\left\|k_{x}\right\|_{\mathcal{H}}^{2}\left\|k_{t}\right\|_{\mathcal{H}}^{2}=K(x, x) K(t, t), x$ and $t \in \mathscr{X}$ (from the CS inequality in $\mathcal{H}$ ), we have $\left|\left\langle g_{\nu} \mid g_{\mu}\right\rangle_{\mathcal{G}}\right| \leqslant \tau_{\nu} \tau_{\mu}$.

Remark 3.5. For $\mu \in \mathcal{T}(K)$ and $\nu \in \mathcal{T}_{\mathbb{C}}(K)$, and denoting by $\nu_{\mathrm{r}}$ and $\nu_{\stackrel{\mathrm{i}}{ }} \in \mathcal{T}(K)$ the real and imaginary parts of $\nu$, we have $g_{\mu}-g_{\nu_{r}} \in \mathcal{G}_{\mathbb{R}}$ and $g_{\nu_{\mathrm{i}}} \in \mathcal{G}_{\mathbb{R}}$. Following Section 2.2, we obtain $\left\|g_{\mu}-g_{\nu}\right\|_{G}^{2}=\left\|g_{\mu}-g_{\nu_{r}}\right\|_{\mathcal{G}}^{2}+\left\|g_{\nu_{\mathrm{i}}}\right\|_{\mathcal{G}}^{2}$, so that $D_{\mu}\left(\nu_{\mathrm{r}}\right) \leqslant D_{\mu}(\nu)$. Consequently, when $L_{\mu}$ is self-adjoint, the search of an approximate measure $\nu$ for the approximation of $L_{\mu}$ by $L_{\nu}$ may be restricted to $\mathcal{T}(K)$.

### 3.3 Nonnegative measures and $L^{2}$-embeddings of RKHSs

Hereafter and in anticipation of the forthcoming Section 4, we further discuss the properties verified by the integral operators considered in Theorem 3.1.

For $\mu \in \mathcal{T}_{+}(K)$, let $L^{2}(\mu)$ be the Hilbert space of all square-integrable functions with respect to $\mu$. From the CS inequality in $\mathcal{H}$, we have

$$
\begin{equation*}
\int_{\mathscr{X}}|h(t)|^{2} \mathrm{~d} \mu(t)=\int_{\mathscr{X}}\left|\left\langle k_{t} \mid h\right\rangle_{\mathcal{H}}\right|^{2} \mathrm{~d} \mu(t) \leqslant\|h\|_{\mathcal{H}}^{2} \tau_{\mu}, h \in \mathcal{H}, \tag{9}
\end{equation*}
$$

so that the linear embedding $l_{\mu}: \mathcal{H} \rightarrow L^{2}(\mu)$, with $l_{\mu}[h]$ the equivalence class of all measurable functions $\mu$-almost everywhere equal to $h$, is bounded (see e.g. [26]).

Lemma 3.3. For all $\mu \in \mathcal{T}_{+}(K)$, the embedding $l_{\mu}$ is HS and $L_{\mu}=l_{\mu}^{*} l_{\mu}$; we in addition have $\operatorname{null}\left(L_{\mu}\right)=\left\{h \in \mathcal{H} \mid l_{\mu}[h]=0\right\}$.

Proof. Let $\left\{h_{i}\right\}_{i \in \mathrm{\rrbracket}}$ be an ONB of $\mathcal{H}$, with $\mathbb{\square} \subseteq \mathbb{N}$. As $K(t, t)=\sum_{i \in \rrbracket}\left|h_{i}(t)\right|^{2}, t \in \mathscr{X}$ (see e.g. [17]), from (6) and by monotone convergence, we have

$$
\sum_{i \in \mathbb{\square}}\left\|l_{\mu}\left[h_{i}\right]\right\|_{L^{2}(\mu)}^{2}=\sum_{i \in \mathbb{\Omega}}\left\langle L_{\mu}\left[h_{i}\right] \mid h_{i}\right\rangle_{\mathcal{H}}=\int_{\mathscr{X}} \sum_{i \in \mathbb{1}}\left|h_{i}(t)\right|^{2} \mathrm{~d} \mu(t)=\tau_{\mu},
$$

so that $l_{\mu}$ is HS. From (6), we also obtain

$$
\left\langle L_{\mu}[h] \mid f\right\rangle_{\mathcal{H}}=\left\langle l_{\mu}[h] \mid l_{\mu}[f]\right\rangle_{L^{2}(\mu)}=\left\langle l_{\mu_{\mu}}^{*} l_{\mu}[h] \mid f\right\rangle_{\mathcal{H}}, h \text { and } f \in \mathcal{H},
$$

and so $L_{\mu}=\iota_{\mu^{\prime}}^{*} l_{\mu}$. Finally, noticing that $\left\langle L_{\mu}[h] \mid h\right\rangle_{\mathcal{H}}=\left\|l_{\mu}[h]\right\|_{L^{2}(\mu)}^{2}$ and that $L_{\mu}$ is PSD, we have $L_{\mu}[h]=0$ if and only if $l_{\mu}[h]=0$.

For $v \in \mathcal{T}_{\mathbb{F}}(K)$, we by definition have $|\nu| \in \mathcal{T}_{+}(K)$ and $\bar{\nu} \in \mathcal{T}_{\mathbb{F}}(K)$; from (6), we can in addition notice that

$$
\left\langle L_{\nu}[h] \mid f\right\rangle_{\mathcal{H}}=\int_{\mathscr{X}} \overline{h(t)} f(t) \mathrm{d} \bar{v}(t)=\left\langle h \mid L_{\bar{v}}[f]\right\rangle_{\mathcal{H}}, h \text { and } f \in \mathcal{H},
$$

so that $L_{v}^{*}=L_{\bar{v}}$. The following relation (Lemma 3.4) holds between the range of $L_{v}$ and the range of $L_{|v|}$.
Lemma 3.4. For all $v \in \mathcal{T}_{\mathbb{F}}(K)$, we have $\overline{\operatorname{range}\left(L_{\nu}\right)}{ }^{\mathcal{H}} \subseteq{\overline{\operatorname{range}\left(L_{|\nu|}\right)}}^{\mathcal{H}}$.
Proof. Using (6) and the CS inequality in $L^{2}(|v|)$, we obtain

$$
\begin{aligned}
\left|\left\langle L_{\bar{v}}[h] \mid f\right\rangle_{\mathcal{H}}\right| & =\left|\int_{\mathscr{X}} \overline{h(t)} f(t) \mathrm{d} v(t)\right| \leqslant \int_{\mathscr{X}}|h(t)||f(t)| \mathrm{d}|\nu|(t) \\
& \leqslant\left\|\iota_{|v|}[h]\right\|_{L^{2}(|v|)}\left\|\iota_{|v|}[f]\right\|_{L^{2}(|v|)}, h \text { and } f \in \mathcal{H} .
\end{aligned}
$$

Lemma 3.3 then entails null $\left(L_{|v|}\right) \subseteq \operatorname{null}\left(L_{\bar{v}}\right)$, and so $\operatorname{null}\left(L_{\bar{v}}\right)^{\perp_{H}} \subseteq \operatorname{null}\left(L_{|\nu|}\right)^{\perp_{H}}$. Since for all $T \in \operatorname{HS}(\mathcal{H})$, we have $\operatorname{null}\left(T^{*}\right)=\operatorname{range}(T)^{\perp_{\mathcal{H}}}$, we conclude by noticing that $L_{\bar{v}}^{*}=L_{v}$ and $L_{|v|}^{*}=L_{|v|}$.

The following Lemma 3.5 illustrates that when $v$ is finitely-supported, then the range of $L_{|v|}$ is fully characterised by the support ${ }^{2}$ of $v$.

Lemma 3.5. For $v=\sum_{i=1}^{n} v_{i} \delta_{s_{i}}$, with $v_{i} \in \mathbb{C}, v_{i} \neq 0$, and $s_{i} \in \mathscr{X}$, we have $\operatorname{range}\left(L_{|v|}\right)=\operatorname{span}_{\mathbb{C}}\left\{k_{s_{1}}, \cdots, k_{s_{n}}\right\}$.

Proof. Noticing that $|v|=\sum_{i=1}^{n}\left|v_{i}\right| \delta_{s_{i}} \in \mathcal{T}_{+}(K)$, from Lemma 3.3, we obtain $\operatorname{null}\left(L_{|v|}\right)=\left\{\left.h \in \mathcal{H}\right|_{|v|}[h]=0\right\}=\bigcap_{i=1}^{n}\left\{h \in \mathcal{H} \mid\left\langle k_{s_{i}} \mid h\right\rangle_{\mathcal{H}}=0\right\}$. Since $L_{|v|}$ is self-adjoint, the result follows.

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## 4 Projections and subspaces defined by measures

Hereafter, we discuss the extent to which the kernel-quadrature setting described in Section 3.2 can be used as a surrogate for the characterisation of closed linear subspaces of $\mathcal{H}$ for the approximation of trace-class integral operators with PSD kernels through projections.

### 4.1 Further notations and general properties

For a closed linear subspace $\mathcal{H}_{S}$ of $\mathcal{H}$, we denote by $P_{S}$ the orthogonal projection from $\mathcal{H}$ onto $\mathcal{H}_{S}$. Endowed with the Hilbert structure of $\mathcal{H}$, the vector space $\mathcal{H}_{S}$ is a RKHS, and its reproducing kernel $K_{S}$ verifies $K_{S}(x, t)=P_{S}\left[k_{t}\right](x), x$ and $t \in \mathscr{X}$.

Remark 4.1. The linear map $T \mapsto P_{S} T$ is the orthogonal projection from $\operatorname{HS}(\mathcal{H})$ onto $\mathcal{R}\left(\mathcal{H}_{S}\right)=\left\{T \in \mathrm{HS}(\mathcal{H}) \mid \operatorname{range}(T) \subseteq \mathcal{H}_{S}\right\}$, the closed linear subspace of $\mathrm{HS}(\mathcal{H})$ of all operators with range included in $\mathcal{H}_{S}$.

Further, the linear map $T \mapsto T P_{S}$ is the orthogonal projection from $\operatorname{HS}(\mathcal{H})$ onto $\mathcal{Z}\left(\mathcal{H}_{S}\right)=\left\{T \in \operatorname{HS}(\mathcal{H}) \mid \operatorname{range}\left(T^{*}\right) \subseteq \mathcal{H}_{S}\right\}$. The two orthogonal projections $T \mapsto P_{S} T$ and $T \mapsto T P_{S}$ commute, and their composition $T \mapsto P_{S} T P_{S}$ is the orthogonal projection from $\operatorname{HS}(\mathcal{H})$ onto $\mathcal{R}\left(\mathcal{H}_{S}\right) \cap \mathcal{Z}\left(\mathcal{H}_{S}\right)$. For all $T \in \operatorname{HS}(\mathcal{H})$, we thus have

$$
\begin{align*}
& \left\|T-P_{S} T\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=\|T\|_{\mathrm{HS}(\mathcal{H})}^{2}-\left\|P_{S} T\right\|_{\mathrm{HS}(\mathcal{H})}^{2}, \text { and }  \tag{10}\\
& \left\|T-P_{S} T P_{S}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=\|T\|_{\mathrm{HS}(\mathcal{H})}^{2}-\left\|P_{S} T P_{S}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}, \tag{11}
\end{align*}
$$

with $\left\|P_{S} T P_{S}\right\|_{\mathrm{HS}(\mathcal{H})} \leqslant\left\|P_{S} T\right\|_{\mathrm{HS}(\mathcal{H})} \leqslant\|T\|_{\mathrm{HS}(\mathcal{H})}$; in particular, if $T$ is self-adjoint, then so is $P_{S} T P_{S}$. Since $\left(P_{S} T\right)^{*}=T^{*} P_{S}$, the orthogonal projections onto $\mathcal{R}\left(\mathcal{H}_{S}\right)$ and $\mathcal{Z}\left(\mathcal{H}_{S}\right)$ are intrinsically related; for this reason, in what follows, we mainly focus on approximations of the form $P_{S} T$ and $P_{S} T P_{S}$.

Lemma 4.1. Let $\mathcal{H}_{S}$ and $\mathcal{H}_{R}$ be two closed linear subspaces of $\mathcal{H}$, with $\mathcal{H}_{R} \subseteq \mathcal{H}_{S}$. For all $T \in \operatorname{HS}(\mathcal{H})$, we have

$$
\left\|P_{R} T\right\|_{\mathrm{HS}(\mathcal{H})} \leqslant\left\|P_{S} T\right\|_{\mathrm{HS}(\mathcal{H})} \quad \text { and } \quad\left\|P_{R} T P_{R}\right\|_{\mathrm{HS}(\mathcal{H})} \leqslant\left\|P_{S} T P_{S}\right\|_{\mathrm{HS}(\mathcal{H})} .
$$

Proof. We denote by $\mathcal{H}_{e}$ the orthogonal complement of $\mathcal{H}_{R}$ in $\mathcal{H}_{S}$. We then have $P_{S}=P_{R}+P_{e}$ and $\left\langle P_{R} T \mid P_{e} T\right\rangle_{\mathrm{HS}(\mathcal{H})}=\left\langle T P_{R} \mid T P_{e}\right\rangle_{\mathrm{HS}(\mathcal{H})}=0, T \in \operatorname{HS}(\mathcal{H})$. It follows that $\left\|P_{S} T\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=\left\|P_{R} T\right\|_{\mathcal{H}}^{2}+\left\|P_{e} T\right\|_{\mathrm{HS}(\mathcal{H})}^{2}$, and that

$$
\begin{aligned}
\left\|P_{S} T P_{S}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=\left\|P_{R} T P_{R}\right\|_{\mathrm{HS}(\mathcal{H})}^{2} & +\left\|P_{e} T P_{e}\right\|_{\mathrm{HS}(\mathcal{H})}^{2} \\
& +\left\|P_{R} T P_{e}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}+\left\|P_{e} T P_{R}\right\|_{\mathrm{HS}(\mathcal{H})}^{2},
\end{aligned}
$$

completing the proof.

Remark 4.2. For a fixed initial operator $T \in \operatorname{HS}(\mathcal{H})$ and a given rank, the subspaces leading to minimal values of $\left\|T-P_{S} T\right\|_{\mathrm{HS}(\mathcal{H})}$ or $\left\|T-P_{S} T P_{S}\right\|_{\mathrm{HS}(\mathcal{H})}$ correspond to spectral truncations of $T$; see e.g. [12, Theorem 4.4.7]. In practical applications, a direct implementation of such optimal approximations is therefore limited to operators for which a singular value decomposition can be computed beforehand.

Orthogonal projections being bounded linear maps, for $\mu \in \mathcal{T}_{\mathbb{F}}(K)$, we have $P_{S} L_{\mu}=\int_{\mathscr{X}} P_{S} S_{k_{t}} \mathrm{~d} \mu(t)$ (see e.g. [28]), that is

$$
P_{S} L_{\mu}[h](x)=\int_{\mathscr{X}} K_{S}(x, t) h(t) \mathrm{d} \mu(t), h \in \mathcal{H} \text { and } x \in \mathscr{X},
$$

and $P_{S} L_{\mu} \in \operatorname{HS}(\mathcal{H})$ may thus be regarded as an integral operator on $\mathcal{H}$ defined by the kernel $K_{S}$ and the measure $\mu$. Since $K_{S}(t, t) \leqslant K(t, t), t \in \mathscr{X}$, we may notice that $\mathcal{T}_{\mathfrak{F}}(K) \subseteq \mathcal{T}_{\mathbb{F}}\left(K_{S}\right)$.

Lemma 4.2. Let $\mathcal{H}_{U}$ and $\mathcal{H}_{V}$ be two closed linear subspaces of $\mathcal{H}$. For all $\mu \in \mathcal{T}_{\mathbb{F}}(K)$, we have

$$
\left\|P_{V} L_{\mu} P_{U}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=\iint_{\mathscr{X}} K_{U}(t, x) K_{V}(x, t) \mathrm{d} \mu(t) \mathrm{d} \bar{\mu}(x) .
$$

Proof. We consider an ONB $\left\{h_{j}\right\}_{j \in \rrbracket}$ of $\mathcal{H}, \square \subseteq \mathbb{N}$. From (6), we have

$$
\begin{align*}
\left\|P_{V} L_{\mu} P_{U}\right\|_{\mathrm{HS}(\mathcal{H})}^{2} & =\sum_{j \in \mathbb{1}}\left\langle L_{\mu} P_{U}\left[h_{i}\right] \mid P_{V} L_{\mu} P_{U}\left[h_{j}\right]\right\rangle_{\mathcal{H}} \\
& =\sum_{j \in \mathbb{1}} \iint_{\mathscr{X}} P_{U}\left[h_{j}\right](t) \overline{P_{U}\left[h_{j}\right](x)} K_{V}(x, t) \mathrm{d} \mu(t) \mathrm{d} \bar{\mu}(x) . \tag{12}
\end{align*}
$$

As $\sum_{j \in 0} P_{U}\left[h_{j}\right](t) \overline{P_{U}\left[h_{j}\right](x)}=K_{U}(t, x), x$ and $t \in \mathscr{X}$ (see e.g. [17]), and since $K_{U}(t, t) \leqslant K(t, t)$ and $K_{V}(t, t) \leqslant K(t, t)$, from the CS inequality in $\ell^{2}(\mathbb{D})$ and in $\mathcal{H}$, we obtain

$$
\iint_{\mathscr{X}} \sum_{j \in \mathbb{1}}\left|P_{U}\left[h_{j}\right](t)\right|\left|P_{U}\left[h_{j}\right](x)\right|\left|K_{V}(x, t)\right| \mathrm{d}|\mu|(t) \mathrm{d}|\mu|(x) \leqslant \tau_{\mu}^{2}
$$

and the result follows form (12) and Fubini's theorem.
Remark 4.3. For $\mu \in \mathcal{T}_{\mathbb{F}}(K)$ and $x \in \mathscr{X}$, we have

$$
\begin{aligned}
& \Gamma\left[P_{S} L_{\mu}\right](x)=\left\langle k_{x} \mid P_{S} L_{\mu}\left[k_{x}\right]\right\rangle_{\mathcal{H}}=\int_{\mathscr{X}} K_{S}(x, t) K(t, x) \mathrm{d} \mu(t), \text { and } \\
& \Gamma\left[P_{S} L_{\mu} P_{S}\right](x)=\left\langle k_{x} \mid P_{S} L_{\mu} P_{S}\left[k_{x}\right]\right\rangle_{\mathcal{H}}=\int_{\mathscr{X}}\left|K_{S}(x, t)\right|^{2} \mathrm{~d} \mu(t)
\end{aligned}
$$

however, for general linear subspaces $\mathcal{H}_{S} \subseteq \mathcal{H}$, the operators $P_{S} L_{\mu}$ and $P_{S} L_{\mu} P_{S}$ do not necessarily belong to $\mathcal{I}(\Gamma)$.

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### 4.2 Approximate measures and projections

Motivated by Lemmas 3.4 and 4.1, for $\nu \in \mathcal{T}_{\mathbb{F}}(K)$, we introduce $\mathcal{H}_{\nu}=\overline{\operatorname{range}\left(L_{|\nu|}\right)}{ }^{\mathcal{H}}$, and we denote by $P_{v}$ the orthogonal projection from $\mathcal{H}$ onto $\mathcal{H}_{v}$. We refer to $\mathcal{H}_{v}$ as the subspace of $\mathcal{H}$ defined by $v$.

Lemma 4.3. For all $v \in \mathcal{T}_{\mathbb{F}}(K)$, we have $L_{v}=P_{v} L_{v}=L_{v} P_{v}=P_{v} L_{v} P_{v}$.
Proof. From Lemma 3.4, we have $L_{v}=P_{v} L_{v}$. Noticing that $L_{\bar{v}}=P_{v} L_{\bar{v}}$, we also have $L_{v}=L_{\bar{v}}^{*}=\left(P_{v} L_{\bar{v}}\right)^{*}=L_{v} P_{v}$, and thus $L_{v}=P_{v} L_{v} P_{v}$.

The following Lemma 4.4 illustrates the relation between the approximations $L_{\nu}, P_{\nu} L_{\mu}$ and $P_{\nu} L_{\mu} P_{\nu}$ of $L_{\mu}$ characterised by an approximate measure $\nu$.

Lemma 4.4. For all $\mu$ and $\nu \in \mathcal{T}_{\mathbb{F}}(K)$, we have

$$
\begin{align*}
\left\|L_{\mu}-L_{v}\right\|_{\mathrm{HS}(\mathcal{H})}^{2} & =\left\|L_{\mu}-P_{v} L_{\mu}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}+\left\|P_{v} L_{\mu}-L_{v}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}, \text { and }  \tag{13}\\
\left\|L_{\mu}-L_{v}\right\|_{\mathrm{HS}(\mathcal{H})}^{2} & =\left\|L_{\mu}-P_{v} L_{\mu} P_{v}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}+\left\|P_{v} L_{\mu} P_{v}-L_{v}\right\|_{\mathrm{HS}(\mathcal{H})}^{2} . \tag{14}
\end{align*}
$$

Proof. Using the notations of Remark 4.1, Lemma 4.3 reads $L_{v} \in \mathcal{R}\left(\mathcal{H}_{\nu}\right) \cap \mathcal{Z}\left(\mathcal{H}_{\nu}\right)$. Noticing that $L_{\mu}-P_{\nu} L_{\mu}$ is orthogonal to $\mathcal{R}\left(\mathcal{H}_{\nu}\right)$ and that $P_{\nu} L_{\mu}-L_{\nu} \in \mathcal{R}\left(\mathcal{H}_{\nu}\right)$, we obtain (13). In the same way, $L_{\mu}-P_{\nu} L_{\mu} P_{\nu}$ is orthogonal to $\mathcal{R}\left(\mathcal{H}_{\nu}\right) \cap \mathcal{Z}\left(\mathcal{H}_{\nu}\right)$, and we have $P_{\nu} L_{\mu} P_{\nu}-L_{\nu} \in \mathcal{R}\left(\mathcal{H}_{\nu}\right) \cap \mathcal{Z}\left(\mathcal{H}_{\nu}\right)$, leading to (14).

For $\mu \in \mathcal{T}_{\mathbb{F}}(K)$, we introduce the two maps $C_{\mu}^{\mathrm{P}}$ and $C_{\mu}^{\mathrm{PP}}: \mathcal{J}_{\mathbb{F}}(K) \rightarrow \mathbb{R}_{\geqslant 0}$, with

$$
C_{\mu}^{\mathrm{P}}(\nu)=\left\|L_{\mu}-P_{v} L_{\mu}\right\|_{\mathrm{HS}(\mathcal{H})}^{2} \text { and } C_{\mu}^{\mathrm{PP}}(v)=\left\|L_{\mu}-P_{v} L_{\mu} P_{\nu}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}, v \in \mathcal{T}_{\mathbb{F}}(K) .
$$

Theorem 4.1. For $\mu \in \mathcal{J}_{\mathbb{F}}(K)$, the map $C_{\mu}^{\mathrm{P}}$ is convex on the real convex cone $\mathcal{T}_{+}(K)$, and for all $v$ and $\eta \in \mathcal{T}_{+}(K)$, we have

$$
\begin{equation*}
\lim _{\rho \rightarrow 0^{+}} \frac{1}{\rho}\left[C_{\mu}^{\mathrm{P}}(v+\rho(\eta-v))-C_{\mu}^{\mathrm{P}}(\nu)\right] \in\{-\infty, 0\} . \tag{15}
\end{equation*}
$$

This statement stays true if we replace $C_{\mu}^{\mathrm{P}}$ by $C_{\mu}^{\mathrm{PP}}$.
Proof. For $v, \eta \in \mathcal{T}_{+}(K)$ and $\rho \in(0,1)$, we set $\xi=v+\rho(\eta-v) \in \mathcal{T}_{+}(K)$. The three operators $L_{v}, L_{\eta}$ and $L_{\xi}$ being PSD, independently of the value of $\rho \in(0,1)$, we have $\operatorname{null}\left(L_{\xi}\right)=\operatorname{null}\left(L_{\nu}\right) \cap \operatorname{null}\left(L_{\eta}\right)$, and so $\mathcal{H}_{\xi}=\overline{\mathcal{H}}_{\nu}+\mathcal{H}_{\eta}{ }^{\mathcal{H}}$. The two maps

$$
\rho \mapsto C_{\mu}^{\mathrm{P}}(v+\rho(\eta-v)) \quad \text { and } \quad \rho \mapsto C_{\mu}^{\mathrm{PP}}(v+\rho(\eta-v))
$$

are therefore constant on the open interval $(0,1)$. From Lemma 4.1 and (10), noticing that $\mathcal{H}_{\nu} \subseteq \mathcal{H}_{\xi}$ and $\mathcal{H}_{\eta} \subseteq \mathcal{H}_{\xi}$, we obtain $C_{\mu}^{\mathrm{P}}(\xi) \leqslant C_{\mu}^{\mathrm{P}}(\nu)$ and $C_{\mu}^{\mathrm{P}}(\xi) \leqslant C_{\mu}^{\mathrm{P}}(\eta)$; from (11), we also get $C_{\mu}^{\mathrm{PP}}(\xi) \leqslant C_{\mu}^{\mathrm{PP}}(\nu)$ and $C_{\mu}^{\mathrm{PP}}(\xi) \leqslant C_{\mu}^{\mathrm{PP}}(\eta)$, concluding the proof.

In view of Theorem 4.1, the maps $C_{\mu}^{\mathrm{P}}$ and $C_{\mu}^{\mathrm{PP}}$ are convex and piecewise-constant on the convex cone $\mathcal{T}_{+}(K)$. In contrast to the map $D_{\mu}$ defined in (7), the directional derivatives of $C_{\mu}^{\mathrm{P}}$ and $C_{\mu}^{\mathrm{PP}}$ are noninformative, in the sense that they do not provide any information on the local steepness of the landscape of these maps; see Figure 1 for an illustration. From Remark 4.1 and Lemma 4.4, we have

$$
\begin{equation*}
C_{\mu}^{\mathrm{P}}(v) \leqslant C_{\mu}^{\mathrm{PP}}(v) \leqslant D_{\mu}(v), v \in \mathcal{T}_{\mathbb{F}}(K) . \tag{16}
\end{equation*}
$$

Also, from a numerical standpoint and in view of Lemma 4.2, for $v \in \mathcal{T}_{\mathbb{F}}(K)$, the evaluation of $C_{\mu}^{\mathrm{P}}(\nu)$ or $C_{\mu}^{\mathrm{PP}}(\nu)$ requires a suitable characterisation of the reproducing kernel $K_{v}$ of $\mathcal{H}_{\nu}$; in practice, $K_{v}$ is a priori unknown and needs to be computed from $K$ and $\nu$. In comparison and in view of (8), the map $D$ merely involves the kernel $K$. See Remark 4.7 for more details.


Figure 1: Graphical representation of the maps $D_{\mu}$ and $C_{\mu}^{\mathrm{PP}}$ as functions of the weight parameters characterising an approximate measure $v \in \mathcal{T}_{+}(K)$. The measures $\mu$ and $\nu$ are supported by the same set of points $\left\{x_{1}, x_{2}\right\} \subseteq \mathscr{X}$, and described by their weight parameters $\left(\omega_{1}, \omega_{2}\right)$ and $\left(v_{1}, v_{2}\right) \in \mathbb{R}_{\geqslant 0}^{2}$, respectively; the red star represents the weight parameters of $\mu=\omega_{1} \delta_{x_{1}}+\omega_{2} \delta_{x_{2}}$. The presented graphs correspond to the case $\omega_{1}=\omega_{2}=1$, with $K$ such that $K\left(x_{1}, x_{1}\right)=1.225, K\left(x_{2}, x_{2}\right)=0.894$ and $K\left(x_{1}, x_{2}\right)=0.316$. In the graph of $C_{\mu}^{\mathrm{PP}}$, the point on the vertical axis indicates the value of the map at $v=0$, and the bold lines indicate the constant values taken by the map along the horizontal axes.

Remark 4.4. For $\mu \in \mathcal{T}_{\mathbb{F}}(K)$ and $s \in \mathscr{X}$, introducing $c_{\delta_{s}}=g_{\mu}(s) /|K(s, s)|^{2}$ if $K(s, s)>0$, and $c_{\delta_{s}}=0$ otherwise, a direct computation yields $P_{\delta_{s}} L_{\mu} P_{\delta_{s}}=c_{\delta_{s}} S_{k_{s}}$. We thus have $P_{\delta_{s}} L_{\mu} P_{\delta_{s}} \in \mathcal{I}(\Gamma)$ and $C_{\mu}^{\mathrm{PP}}\left(\delta_{s}\right)=D_{\mu}\left(c_{\delta_{s}} \delta_{s}\right)$. Consequently, in Figure 1, the graph of $C_{\mu}^{\mathrm{PP}}$ is tangent to the graph of $D_{\mu}$ along the horizontal axes.

### 4.3 Partial $L^{2}$-embeddings

Following Section 3.3, for $\mu \in \mathcal{T}_{+}(K)$, the embedding ${ }_{\mu}: \mathcal{H} \rightarrow L^{2}(\mu)$ is HS. For $f \in L^{2}(\mu)$ and $x \in \mathscr{X}$, we have

$$
\left\langle k_{x} \mid l_{\mu}^{*}[f]\right\rangle_{\mathcal{H}}=\left\langle l_{\mu}\left[k_{x}\right] \mid f\right\rangle_{L^{2}(\mu)}=\int_{\mathscr{X}} K(x, t) f(t) \mathrm{d} \mu(t) .
$$

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Consequently, in addition to $L_{\mu}={l_{\mu}^{*}}_{\mu} l_{\mu} \in \operatorname{HS}(\mathcal{H})$, the three operators

$$
\begin{equation*}
l_{\mu}^{*}: L^{2}(\mu) \rightarrow \mathcal{H}, \quad l_{\mu} \mu_{\mu}^{*}: L^{2}(\mu) \rightarrow L^{2}(\mu), \quad \text { and } \quad l_{\mu} \iota_{\mu_{\mu}}^{*} l_{\mu}: \mathcal{H} \rightarrow L^{2}(\mu) \tag{17}
\end{equation*}
$$

can also be regarded as integral operators defined by the kernel $K$ and the nonnegative measure $\mu$. These four interpretations are inherent to $K$, which characterises $\mathcal{H}$, and $\mu$, which characterises $L^{2}(\mu)$; see for instance $[4,22,19,26,20]$ for illustrations. In each case, the corresponding operator is HS , and we denote by $\operatorname{HS}(\mu, \mathcal{H}), \mathrm{HS}(\mu)$ and $\operatorname{HS}(\mathcal{H}, \mu)$ the Hilbert spaces of all HS operators from $L^{2}(\mu)$ to $\mathcal{H}$, on $L^{2}(\mu)$, and from $\mathcal{H}$ to $L^{2}(\mu)$, respectively.

For a closed linear subspace $\mathcal{H}_{S} \subseteq \mathcal{H}$, the embedding $l_{\mu}$ can be approximated by the partial embedding $l_{\mu} P_{S}$; for $f \in L^{2}(\mu)$ and $x \in \mathscr{X}$, we in particular have

$$
\left\langle k_{x} \mid P_{S} l_{\mu}^{*}[f]\right\rangle_{\mathcal{H}}=\left\langle l_{\mu} P_{S}\left[k_{x}\right] \mid f\right\rangle_{L^{2}(\mu)}=\int_{\mathscr{X}} K_{S}(x, t) f(t) \mathrm{d} \mu(t) .
$$

For $L_{\mu}=l_{\mu}^{*} l_{\mu} \in \operatorname{HS}(\mathcal{H})$, we obtain the approximation $P_{S} l_{\mu}^{*} l_{\mu} P_{S}$ discussed in Section 4.2, and for the three operators defined in (17), we obtain the approximations

$$
P_{S} l_{\mu}^{*} \in \operatorname{HS}(\mu, \mathcal{H}), \quad l_{\mu} P_{S} l_{\mu}^{*} \in \operatorname{HS}(\mu), \quad \text { and } \quad l_{\mu} P_{S} l_{\mu}^{*} l_{\mu} P_{S} \in \operatorname{HS}(\mathcal{H}, \mu) .
$$

In what follows and for simplicity, we mainly focus on the approximations related to $\operatorname{HS}(\mu, \mathcal{H})$ and $\operatorname{HS}(\mu)$; the case of $\operatorname{HS}(\mathcal{H}, \mu)$ is more briefly discussed in Remark 4.5.

Lemma 4.5. Let $\mathcal{H}_{S}$ be a closed linear subspace of $\mathcal{H}$. For $\mu \in \mathcal{T}_{+}(K)$, we have

$$
\begin{align*}
& \left\|l_{\mu}^{*}-P_{S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu, \mathcal{H})}^{2}=\int_{\mathscr{X}} K(t, t)-K_{S}(t, t) \mathrm{d} \mu(t), \text { and }  \tag{18}\\
& \left\|l_{\mu} \mu_{\mu}^{*}-l_{\mu} P_{S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)}^{2}=\iint_{\mathscr{X}}\left|K(x, t)-K_{S}(x, t)\right|^{2} \mathrm{~d} \mu(t) \mathrm{d} \mu(x) . \tag{19}
\end{align*}
$$

Proof. Let $\mathcal{H}_{0 S}$ be the orthogonal complement of $\mathcal{H}_{S}$ in $\mathcal{H}$; endowed with the Hilbert structure of $\mathcal{H}, \mathcal{H}_{0 S}$ is a RKHS, and $K_{0 S}=K-K_{S}$. For an ONB $\left\{h_{j}\right\}_{j \in \emptyset}$ of $\mathcal{H}$, $\llbracket \subseteq \mathbb{N}$, we have

$$
\left\|l_{\mu}^{*}-P_{S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu, \mathcal{H})}^{2}=\left\|l_{\mu} P_{0 S}\right\|_{\mathrm{HS}(\mathcal{H}, \mu)}^{2}=\sum_{j \in \mathrm{\square}} \int_{\mathscr{X}}\left|P_{0 S}\left[h_{j}\right](t)\right|^{2} \mathrm{~d} \mu(t) ;
$$

since $\sum_{j \in \mathbb{Q}}\left|P_{0 S}\left[h_{j}\right](t)\right|^{2}=K_{0 S}(t, t), t \in \mathscr{X}$, equality (18) follows by monotone convergence. We also have

$$
\left\|l_{\mu} l_{\mu}^{*}-l_{\mu} P_{S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)}^{2}=\left\|l_{\mu} P_{0 S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)}^{2}=\left\|P_{0 S} L_{\mu} P_{0 S}\right\|_{\mathrm{HS}(\mathcal{H})}^{2},
$$

so that (19) follows from Lemma 4.2.
The following inequality (Lemma 4.6) holds between the approximations in $\operatorname{HS}(\mu)$ and $\operatorname{HS}(\mathcal{H})$ defined by a subspace $\mathcal{H}_{S}$. We recall that $\left\|l_{\mu} \mu_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)}=\left\|l_{\mu^{*}}^{*}{ }_{\mu}\right\|_{\mathrm{HS}(\mathcal{H})}$, and that $\left\|l_{\mu}^{*} l_{\mu}-P_{S} l_{\mu}^{*} l_{\mu}\right\|_{\mathrm{HS}(\mathcal{H})} \leqslant\left\|l_{\mu}^{*} l_{\mu}-P_{S} l_{\mu}^{*} l_{\mu} P_{S}\right\|_{\mathrm{HS}(\mathcal{H})}$; see Remark 4.1.

Lemma 4.6. Let $\mathcal{H}_{S}$ be a closed linear subspace of $\mathcal{H}$. For all $\mu \in \mathcal{T}_{+}(K)$, we have

$$
\left\|t_{\mu} l_{\mu}^{*}-l_{\mu} P_{S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)} \leqslant\left\|l_{\mu}^{*} l_{\mu}-P_{S}{ }_{l_{\mu}}^{*} l_{\mu}\right\|_{\mathrm{HS}(\mathcal{H})} .
$$

Proof. We have $\left\|l_{\mu} l_{\mu}^{*}-l_{\mu} P_{S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)}=\left\|P_{0 S} L_{\mu} P_{0 S}\right\|_{\mathrm{HS}(\mathcal{H})}$, with $P_{0 S}=\mathrm{id}_{\mathcal{H}}-P_{S}$. The operators $P_{0 S} L_{\mu} P_{0 S}$ and $P_{0 S} L_{\mu} P_{S}$ being orthogonal in $\operatorname{HS}(\mathcal{H})$, we get

$$
\left\|P_{0 S} L_{\mu}\left(P_{0 S}+P_{S}\right)\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=\left\|P_{0 S} L_{\mu} P_{0 S}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}+\left\|P_{0 S} L_{\mu} P_{S}\right\|_{\mathrm{HS}(\mathcal{H})}^{2},
$$

so that $\left\|l_{\mu} l_{\mu}^{*}-l_{\mu} P_{S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)} \leqslant\left\|l_{\mu}^{*} l_{\mu}-P_{S} l_{\mu}^{*} l_{\mu}\right\|_{\mathrm{HS}(\mathcal{H})}$.
Lemma 4.7. Let $\mathcal{H}_{S}$ and $\mathcal{H}_{R}$ be two closed linear subspaces of $\mathcal{H}$, with $\mathcal{H}_{R} \subseteq \mathcal{H}_{S}$. For all $\mu \in \mathcal{T}_{+}(K)$, we have
$\left\|l_{\mu}^{*}-P_{S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu, \mathcal{H})} \leqslant\left\|l_{\mu}^{*}-P_{R} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu, \mathcal{H})}$ and $\left\|l_{l^{\prime} l_{\mu}^{*}-l_{\mu}} P_{S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)} \leqslant\left\|l_{\mu} l_{\mu}^{*}-l_{\mu} P_{R} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)}$.
Proof. We denote by $\mathcal{H}_{e}$ the orthogonal complement of $\mathcal{H}_{R}$ in $\mathcal{H}_{S}$. Noticing that $\left\langle l_{\mu}^{*}-P_{R} l_{\mu}^{*} \mid P_{e} e_{\mu}^{*}\right\rangle_{\mathrm{HS}(\mu, \mathcal{H})}=\left\|P_{e} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu, \mathcal{H})}^{2}$, we obtain

$$
\left\|l_{\mu}^{*}-P_{S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu, \mathcal{H})}^{2}=\left\|l_{\mu}^{*}-P_{R} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu, \mathcal{H})}^{2}-\left\|P_{e} e_{\mu}^{*}\right\|_{\mathrm{HS}(\mu, \mathcal{H})}^{2} .
$$

Denoting by $\mathcal{H}_{0 S}$ and $\mathcal{H}_{0 R}$ the orthogonal complements of $\mathcal{H}_{S}$ and $\mathcal{H}_{R}$ in $\mathcal{H}$, respectively, we have $\mathcal{H}_{0 S} \subseteq \mathcal{H}_{0 R}$, and Lemma 4.1 gives

$$
\left\|l_{\mu} P_{0 S} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)}=\left\|P_{0 S} L_{\mu} P_{0 S}\right\|_{\mathrm{HS}(\mathcal{H})} \leqslant\left\|P_{0 R} L_{\mu} P_{0 R}\right\|_{\mathrm{HS}(\mathcal{H})}=\left\|l_{\mu} P_{0 R}{ }^{*}{ }_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)},
$$

which completes the proof.
Following Section 4.2 and considering subspaces of $\mathcal{H}$ defined by measures, we introduce the maps $C_{\mu}^{\mathrm{tr}}$ and $C_{\mu}^{\mathrm{F}}: \mathcal{T}_{\mathbb{F}}(K) \rightarrow \mathbb{R}_{\geqslant 0}$, with

$$
C_{\mu}^{\operatorname{tr}}(\nu)=\left\|l_{\mu}^{*}-P_{\nu} l_{\mu}^{*}\right\|_{\mathrm{HS}(\mu, \mathcal{H})}^{2} \text { and } C_{\mu}^{\mathrm{F}}(\nu)=\left\|\iota_{\mu} \iota_{\mu}^{*}-\iota_{\mu} P_{\nu} v_{\mu}^{*}\right\|_{\mathrm{HS}(\mu)}^{2}, \nu \in \mathcal{T}_{\mathbb{F}}(K) .
$$

The notations $C_{\mu}^{\mathrm{tr}}$ and $C_{\mu}^{\mathrm{F}}$ are motivated by the relation existing between these maps and the trace and Frobenius norms for matrices; see Section 4.4.
Corollary 4.1. The statement of Theorem 4.1 also holds for the maps $C_{\mu}^{\mathrm{tr}}$ and $C_{\mu}^{\mathrm{F}}$, that is, these two maps are convex on the convex cone $\mathcal{T}_{+}(K)$, and their directional derivatives take values in the set $\{-\infty, 0\}$.

Proof. Following the proof of Theorem 4.1, for two measures $v$ and $\eta \in \mathcal{T}_{+}(K)$ and for $\rho \in(0,1)$, we set $\xi=v+\rho(\eta-v) \in \mathcal{T}_{+}(K)$. We then have $\mathcal{H}_{\xi}=\overline{\mathcal{H}}_{\nu}+\mathcal{H}_{\eta}{ }^{\mathcal{H}}$, and the result follows from Lemma 4.7.

Remark 4.5. We consider four closed linear subspaces $\mathcal{H}_{R}, \mathcal{H}_{S}, \mathcal{H}_{U}$ and $\mathcal{H}_{V}$ of $\mathcal{H}$, and let $\left\{h_{j}\right\}_{j \in 0}$ be an ONB of $\mathcal{H}$. For $\mu \in \mathcal{T}_{+}(K)$, we have

$$
\begin{aligned}
& \left\langle l_{\mu} P_{R} l_{\mu}^{*} l_{\mu} P_{S} \mid l_{\mu} P_{U} l_{\mu}^{*} l_{\mu} P_{V}\right\rangle_{\mathrm{HS}(\mathcal{H}, \mu)} \\
& \quad=\sum_{j \in \mathrm{Q}} \iiint \int_{\mathscr{X}} P_{V}\left[h_{j}\right](t) \overline{P_{S}\left[h_{j}\right](s)} K_{R}(s, x) K_{U}(x, t) \mathrm{d} \mu(s) \mathrm{d} \mu(t) \mathrm{d} \mu(x),
\end{aligned}
$$

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with, form the CS inequality in $\ell^{2}(\mathbb{D})$ and in $\mathcal{H}$,

$$
\iiint_{\mathscr{X}} \sum_{j \in \mathbb{I}}\left|P_{V}\left[h_{j}\right](t)\right|\left|P_{S}\left[h_{j}\right](s)\right|\left|K_{R}(s, x)\right|\left|K_{U}(x, t)\right| \mathrm{d}|\mu|(s) \mathrm{d}|\mu|(t) \mathrm{d}|\mu|(x) \leqslant \tau_{\mu}^{3}
$$

The approximation error $\left\|l_{\mu} l_{\mu}^{*} l_{\mu}-l_{\mu} P_{S} l_{\mu}^{*} l_{\mu} P_{S}\right\|_{\mathrm{HS}(\mathcal{H}, \mu)}^{2}$ may hence be expressed as a triple integral involving the kernels $K$ and $K_{S}$.

Through the partial embedding $l_{\mu} P_{\nu}$, a measure $\nu \in \mathcal{J}_{\mathbb{F}}(K)$ defines approximations of the integral operators (with kernel $K$ and measure $\mu$ ) induced by $l_{\mu}$ in the four Hilbert spaces $\operatorname{HS}(\mathcal{H}), \operatorname{HS}(\mu, \mathcal{H}), \operatorname{HS}(\mu)$ and $\operatorname{HS}(\mathcal{H}, \mu) ;$ since these approximations are all composite of $t_{\mu} P_{v}$, following Section 4.2, the map $D_{\mu}$ may more generally be used as a surrogate for the characterisation of low-rank approximations of these operators. From Lemma 4.6, we have $C_{\mu}^{\mathrm{F}}(\nu) \leqslant C_{\mu}^{\mathrm{P}}(\nu), v \in \mathcal{J}_{\mathbb{F}}(K)$, adding to the sequence of inequalities (16). In addition to the informative nature of the directional derivatives of $D_{\mu}$, the exploration of the kernel-quadrature approximation landscape does not require the characterisation of orthogonal projections onto subspaces of $\mathcal{H}$; this feature is of interest from numerical-complexity standpoint, as illustrated in the following Remarks 4.6 and 4.7.

Remark 4.6. In the framework of Lemma 3.5, for a discrete measure $v$ supported by a set of $n \in \mathbb{N}$ points $S=\left\{s_{1}, \cdots, s_{n}\right\} \subseteq \mathscr{X}$, the reproducing kernel $K_{v}$ of $\mathcal{H}_{v}$ can be written as

$$
K_{\downarrow}(x, t)=\sum_{i, j=1}^{n} K\left(x, s_{i}\right) \varkappa_{i, j} K\left(s_{j}, t\right), x \text { and } t \in \mathscr{X},
$$

where $\varkappa_{i, j}$ is the $i, j$ entry of the pseudoinverse of the $n \times n$ kernel matrix with $i, j$ entry $K\left(s_{i}, s_{j}\right)$; see (20) in Section 4.4 for an illustration. The numerical complexity of the evaluation of $K_{\nu}$ at $M \in \mathbb{N}$ distinct locations in $\mathscr{X} \times \mathscr{X}$ is thus $\mathcal{O}\left(n^{3}+n^{2} M\right)$. The term $\mathcal{O}\left(n^{3}\right)$ is related to the pseudoinversion of the kernel matrix defined by $K$ and $S$, while the term $\mathcal{O}\left(n^{2} M\right)$ corresponds to the evaluation, from this pseudoinverse and the kernel $K$, of $K_{v}$ at $M$ different locations.

Remark 4.7. Assuming that $\mu$ and $\nu$ are supported by $N$ and $n$ points, respectively, from (8), the computational complexity of the evaluation of $D_{\mu}(\nu)$ up to the constant $\left\|g_{\mu}\right\|_{\mathcal{G}}^{2}$ is $\mathcal{O}\left(n^{2}+n N\right)$. In comparison, following Remark 4.6 and assuming that the kernel $K_{v}$ is evaluated from $K$ and $v$, the complexity is $\mathcal{O}\left(n^{3}+n^{2} N\right)$ for $C_{\mu}^{\text {tr }}(\nu)$, and $\mathcal{O}\left(n^{3}+n^{2} N^{2}\right)$ for $C_{\mu}^{\mathrm{P}}(\nu), C_{\mu}^{\mathrm{PP}}(\nu)$ and $C_{\mu}^{\mathrm{F}}(\nu)$.

If we for instance in addition assume that $\mathscr{X}$ is a discrete set of size $N$ (and so, that $n \leqslant N$; see Section 4.4 for an illustration), then the numerical cost of multiple evaluations of $D_{\mu}$ may be reduced by beforehand computing and storing the values $\left\{g_{\mu}(x)\right\}_{x \in \mathscr{X}}$ of the target potential $g_{\mu} \in \mathcal{C}$; this operation has complexity $\mathcal{O}\left(N^{2}\right)$, but can nevertheless be easily parallelised. From (8), the evaluation of $\left\langle g_{\nu} \mid g_{\mu}\right\rangle_{G}$ then reduces to the integration of $g_{\mu}$ with respect to $\overline{\bar{v}}$, with complexity $\mathcal{O}(n)$.

### 4.4 Column sampling for PSD-matrix approximation

Let $\mathbf{K}$ be a $N \times N$ PSD matrix, with $N \in \mathbb{N}$; we denote by $[N]$ the set of all integers between 1 and $N$. For a subset $I \subseteq[N]$ the Nyström approximation ${ }^{3}$ of $\mathbf{K}$ induced by the columns of $\mathbf{K}$ with index in $I$ is the $N \times N$ PSD matrix

$$
\begin{equation*}
\hat{\mathbf{K}}(I)=\mathbf{K}_{\cdot, I}\left(\mathbf{K}_{I, I}\right)^{\dagger} \mathbf{K}_{I, 0}, \tag{20}
\end{equation*}
$$

where $\mathbf{K}_{, I}$ is the $N \times n$ matrix defined by the columns of $\mathbf{K}$ with index in $I$, and where $\left(\mathbf{K}_{I, I}\right)^{\dagger}$ is the pseudoinverse of the principal submatrix of $\mathbf{K}$ defined by $I$ (and $\mathbf{K}_{I}$. consists of rows of $\mathbf{K}$; it is the conjugate-transpose of $\mathbf{K}_{\cdot, I}$ ); see e.g. [9,18, 14, 10].

For $i$ and $j \in[N]$, the $i, j$ entry of $\mathbf{K}$ may be regarded as the value $K(i, j)$ of a PSD kernel $K$ defined on the discrete set $\mathscr{X}=[N]$. The $j$-th column of $\mathbf{K}$ then corresponds to the function $k_{j} \in \mathcal{H}$, with $j \in \mathscr{X}$, and the subset $I$ defines the linear subspace $\mathcal{H}_{I}=\operatorname{span}_{\mathbb{C}}\left\{k_{j} \mid j \in I\right\} \subseteq \mathcal{H}$; in particular, the $i, j$ entry of $\hat{\mathbf{K}}(I)$ is $K_{I}(i, j)$, with $K_{I}$ the reproducing kernel of $\mathcal{H}_{I}$.

Introducing $\mu=\sum_{i=1}^{N} \delta_{i}$, the Hilbert space $L^{2}(\mu)$ can be identified with the Euclidean space $\mathbb{C}^{N}$; following Section 4.3, we can then notice that

- the trace norm $\|\mathbf{K}-\hat{\mathbf{K}}(I)\|_{\text {tr }}$ corresponds to (18), and
- the squared Frobenius norm $\|\mathbf{K}-\hat{\mathbf{K}}(I)\|_{\mathrm{F}}^{2}$ corresponds to (19).

The column-sampling problem for the Nyström approximation of a PSD matrix $\mathbf{K}$ (that is, the search of a subset $I \subseteq[N]$ leading to an efficient approximation $\hat{\mathbf{K}}(I)$ of $\mathbf{K}$ ) is thus a special instance of the general framework discussed in Section 4.3. In particular, the support of an approximate measure $v$ on $\mathscr{X}=[N]$ defines a subset of columns of $\mathbf{K}$, and the kernel-quadrature setting may be used as surrogate for the characterisation of such measures.

## 5 Concluding discussion

We described the overall framework surrounding the isometric representation of trace-class integral operators with PSD kernel $K$ as potentials in the RKHS $\mathcal{G}$ associated with the kernel $|K|^{2}$, and illustrated the parallel between the quadrature approximation of such operators and kernel quadrature rules in $\mathcal{G}$. Through subspaces defined by measures and partial $L^{2}$-embeddings, we discussed the extent to which the kernel-quadrature setting can be used as a surrogate to characterise projection-based approximations of general integral operators with PSD kernels. We also illustrated the connections between the considered framework and the low-rank approximation of PSD matrices through column sampling.

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In contrast to the projection-based error maps $C_{\mu}^{\mathrm{P}}, C_{\mu}^{\mathrm{PP}}, C_{\mu}^{\mathrm{tr}}$ and $C_{\mu}^{\mathrm{F}}$ considered in Section 4, the quadrature-based error map $D_{\mu}$ defined in (7) is convex and admits informative directional derivatives on any convex set of approximate measures; to this extent, the map $D_{\mu}$ may be regarded as a differentiable relaxation of the piecewise-constant projection-based error maps of Section 4. Further, from a numerical standpoint and as opposed to the projection setting, the exploration of the kernelquadrature approximation landscape does not require the characterisation of orthogonal projections onto subspaces of $\mathcal{H}$, and instead only involves the evaluation of integrals related to $|K|^{2}$. All together, these properties make the kernel-quadrature setting an interesting framework for the design of numerically-efficient problemdependent sampling strategies for the low-rank approximation of trace-class integral operators with PSD kernels.

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[^1]:    ${ }^{1}$ We only consider finite complex measures, while signed measures may not necessarily be finite.

[^2]:    ${ }^{2}$ Assuming that $\mathscr{X}$ is a topological space and that the RKHS $\mathcal{H}$ consists of continuous functions, a similar characterisation could be obtained for Radon measures, for instance. In this note, we merely focus on the finitely-supported case as it does not require any further assumption and is one of the most relevant in terms of practical applications.

[^3]:    ${ }^{3}$ In the machine-learning literature, Nyström approximation refers to the low-rank approximation of PSD matrices through column sampling; although related, this terminology should not to be confused with the quadrature method for the numerical approximation of integral equations.

