# Approximate Log-Hilbert-Schmidt Distances between Covariance Operators for Image Classification Supplementary Material 

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#### Abstract

The Supplementary Material contains the proofs for all mathematical results stated in the main paper. We also describe in more detail the Quasi-random Fourier feature map approach in Section 3.1 of the main paper.


## 1. Proofs for main mathematical results

For clarity, we restate all the mathematical results that we wish to prove here.
Theorem 1. Assume that $\gamma \neq \mu, \gamma>0, \mu>0$. Then

$$
\begin{equation*}
\lim _{D \rightarrow \infty}\left\|\log \left(C_{\hat{\Phi}_{D}(\mathbf{x})}+\gamma I_{D}\right)-\log \left(C_{\hat{\Phi}_{D}(\mathbf{y})}+\mu I_{D}\right)\right\|_{F}=\infty \tag{1}
\end{equation*}
$$

Theorem 2. Assume that $\gamma=\mu>0$ and that $\lim _{D \rightarrow \infty} \hat{K}_{D}(x, y)=K(x, y)$ for every pair $(x, y) \in \mathcal{X} \times \mathcal{X}$. Then

$$
\begin{equation*}
\lim _{D \rightarrow \infty}\left\|\log \left(C_{\hat{\Phi}_{D}(\mathbf{x})}+\gamma I_{D}\right)-\log \left(C_{\hat{\Phi}_{D}(\mathbf{y})}+\gamma I_{D}\right)\right\|_{F}=\left\|\log \left(C_{\Phi(\mathbf{x})}+\gamma I_{\mathcal{H}}\right)-\log \left(C_{\Phi(\mathbf{y})}+\gamma I_{\mathcal{H}}\right)\right\|_{\mathrm{eHS}} \tag{2}
\end{equation*}
$$

We need the following preliminary results. Let $\mathcal{H}$ be a separable Hilbert space, equipped with norm $\|\|$, and $A: \mathcal{H} \rightarrow \mathcal{H}$ be a bounded linear operator. We recall that the operator norm of $A$ is defined to be

$$
\begin{equation*}
\|A\|=\sup _{x \neq 0} \frac{\|A x\|}{\|x\|} . \tag{3}
\end{equation*}
$$

If $A$ is self-adjoint, compact, and positive, then

$$
\begin{equation*}
\|A\|=\lambda_{\max }(A) \tag{4}
\end{equation*}
$$

where $\lambda_{\max }(A)$ denotes the largest eigenvalue of $A$. The trace norm of $A$ in this case is given by

$$
\begin{equation*}
\|A\|_{\operatorname{tr}}=\sum_{k=1}^{\infty} \lambda_{k}(A)=\operatorname{tr}(A) \tag{5}
\end{equation*}
$$

Lemma 1. Let $\mathcal{H}$ be a separable Hilbert space. Let $r \in \mathbb{N}$ be fixed. Let $A \in \mathcal{L}(\mathcal{H})$ be a self-adjoint, positive operator with finite rank $r<\infty$. Then

$$
\begin{array}{r}
\|A\| \leq\|A\|_{\mathrm{HS}} \leq \sqrt{r}\|A\| \\
\|A\| \leq\|A\|_{\mathrm{tr}} \leq r\|A\| . \tag{7}
\end{array}
$$

Thus convergences in the $\left\|\|_{\mathrm{HS}}\right.$ norm, the $\| \|_{\mathrm{tr}}$ norm, and the $\|\|$ norm are all equivalent to each other.

Proof. By definition of the || || and || || ${ }_{\text {HS }}$ norms and the finite rank assumption, we have

$$
\|A\|^{2}=\lambda_{\max }^{2}(A) \leq \sum_{j=1}^{r} \lambda_{j}^{2}(A)=\|A\|_{\mathrm{HS}}^{2} \leq r \lambda_{\max }^{2}(A)
$$

from which the first inequality follows. Similarly, for the second inequality, we have

$$
\|A\|=\lambda_{\max }(A) \leq \sum_{j=1}^{r} \lambda_{k}(A)=\|A\|_{\operatorname{tr}} \leq r \lambda_{\max }(A)=r\|A\| .
$$

This completes the proof of the lemma.
Lemma 2. Let $\mathcal{H}$ be a separable Hilbert space. Let $r \in \mathbb{N}$ be fixed. Let $A,\left\{A_{k}\right\}_{k \in \mathbb{N}}$ be self-adjoint, positive operators of rank at most $r$, such that $\lim _{k \rightarrow \infty}\left\|A_{k}-A\right\|_{\mathrm{HS}}=0$. Then

$$
\begin{equation*}
\lim _{k \rightarrow \infty}\left\|\log \left(I+A_{k}\right)-\log (I+A)\right\|_{\mathrm{HS}}=0 \tag{8}
\end{equation*}
$$

Proof. By assumption, the operators in the sequence $\left(A_{k}-A\right)_{k \in \mathbb{N}}$ all have rank at most $2 r$. Thus from Lemma 1 , the convergence $\left\|A_{k}-A\right\|_{\text {HS }}$ is equivalent to the convergence $\left\|A_{k}-A\right\|$.

The operators in the sequence $\left(\log \left(I+A_{k}\right)\right)_{k \in \mathbb{N}}$ are also self-adjoint, positive, and of rank at most $r$. The operators in the sequence $\left(\log \left(I+A_{k}\right)-\log (I+A)\right)_{k \in \mathbb{N}}$ have rank at most $2 r$ and thus the convergence $\left\|\log \left(I+A_{k}\right)-\log (I+A)\right\|_{\text {HS }}$ is equivalent to the convergence $\left\|\log \left(I+A_{k}\right)-\log (I+A)\right\|$. Thus we have

$$
\left\|A_{k}-A\right\|_{\mathrm{HS}} \rightarrow 0 \Longleftrightarrow\left\|A_{k}-A\right\| \rightarrow 0 \Longleftrightarrow \lambda_{\max }\left(A_{k}\right) \rightarrow \lambda_{\max }(A)
$$

It follows that

$$
\begin{aligned}
\log \left(1+\lambda_{\max }\left(A_{k}\right)\right) \rightarrow \log \left(1+\lambda_{\max }(A)\right) & \Longleftrightarrow\left\|\log \left(I+A_{k}\right)-\log (I+A)\right\| \rightarrow 0 \\
& \Longleftrightarrow\left\|\log \left(I+A_{k}\right)-\log (I+A)\right\|_{\mathrm{HS}} \rightarrow 0
\end{aligned}
$$

This completes the proof of the lemma.
Lemma 3. Let $\mathcal{H}$ be a separable Hilbert space. Let $r \in \mathbb{N}$ be fixed. Let $A,\left\{A_{k}\right\}_{k \in \mathbb{N}}$ be self-adjoint, positive operators of rank at most $r$, such that $\lim _{k \rightarrow \infty}\left\|A_{k}-A\right\|_{\mathrm{HS}}=0$. Then

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \operatorname{tr}\left[\log \left(I+A_{k}\right)-\log (I+A)\right]=0 \tag{9}
\end{equation*}
$$

Proof. By assumption, the operators in the sequence $\left(\log \left(I+A_{k}\right)\right)_{k \in \mathbb{N}}$ are also self-adjoint, positive, and of rank at most $r$. The operators in the sequence $\left(\log \left(I+A_{k}\right)-\log (I+A)\right)_{k \in \mathbb{N}}$ have rank at most $2 r$ and by Lemma $2, \lim _{k \rightarrow \infty} \| \log (I+$ $\left.A_{k}\right)-\log (I+A) \|_{\mathrm{HS}}=0$. By Lemma 1, this convergence is equivalent to convergence in the $\left\|\left\|\|_{\mathrm{tr}}\right.\right.$ norm. Thus we have

$$
\left|\operatorname{tr}\left[\log \left(I+A_{k}\right)-\log (I+A)\right]\right| \leq\left\|\log \left(I+A_{k}\right)-\log (I+A)\right\|_{\mathrm{tr}} \rightarrow 0
$$

as $k \rightarrow \infty$. This completes the proof of the lemma.
Lemma 4. Let $\mathcal{H}$ be a separable Hilbert space. Let $r \in \mathbb{N}$ be fixed. Let $A,\left\{A_{k}\right\}_{k \in \mathbb{N}}, B,\left\{B_{k}\right\}_{k \in \mathbb{N}}$ be self-adjoint, positive operators of rank at most $r$, such that $\lim _{k \rightarrow \infty}\left\|A_{k}-A\right\|_{\mathrm{HS}}=0$ and $\lim _{k \rightarrow \infty}\left\|B_{k}-B\right\|_{\mathrm{HS}}=0$. Then

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \operatorname{tr}\left[\log \left(I+A_{k}\right) \log \left(I+B_{k}\right)\right]=\operatorname{tr}[\log (I+A) \log (I+B)] \tag{10}
\end{equation*}
$$

Proof. From Lemma 2, we have

$$
\begin{aligned}
& \lim _{k \rightarrow \infty}\left\|\log \left(I+A_{k}\right)-\log (I+A)\right\|_{\mathrm{HS}}=0 \\
& \lim _{k \rightarrow \infty}\left\|\log \left(I+B_{k}\right)-\log (I+B)\right\|_{\mathrm{HS}}=0
\end{aligned}
$$

Thus, using Cauchy-Schwarz inequality and the definition $\langle A, B\rangle_{\mathrm{HS}}=\operatorname{tr}\left(A^{T} B\right)$, we obtain

Taking limit on both sides as $k \rightarrow \infty$, we obtain

$$
\lim _{k \rightarrow \infty} \operatorname{tr}\left[\log \left(I+A_{k}\right) \log \left(I+B_{k}\right)-\log (I+A) \log (I+B)\right]=0
$$

This completes the proof of the lemma.
Lemma 5. [2] Let $\mathcal{H}_{1}$ and $\mathcal{H}_{2}$ be two separable Hilbert spaces. Let $A: \mathcal{H}_{1} \rightarrow \mathcal{H}_{2}$ and $B: \mathcal{H}_{2} \rightarrow \mathcal{H}_{1}$ be two bounded linear operators. Then the nonzero eigenvalues of $B A: \mathcal{H}_{1} \rightarrow \mathcal{H}_{1}$ and $A B: \mathcal{H}_{2} \rightarrow \mathcal{H}_{2}$, if they exist, are the same.

In the following, we identify $\mathcal{H}$ with $\ell^{2}$ and $\mathbb{R}^{D}$ with a $D$-dimensional subspace of $\ell^{2}$, that is

$$
\begin{equation*}
a=\left(a_{j}\right)_{j=1}^{D} \in \mathbb{R}^{D} \Longleftrightarrow a=\left(a_{1}, \ldots, a_{D}, 0,0, \ldots\right) \in \ell^{2} . \tag{11}
\end{equation*}
$$

For the data matrices $\mathbf{x}=\left[x_{1}, \ldots, x_{m}\right], \mathbf{y}=\left[y_{1}, \ldots, y_{m}\right]$, let $K[\mathbf{x}], K[\mathbf{y}], \hat{K}_{D}[\mathbf{x}], \hat{K}_{D}[\mathbf{y}]$ be the $m \times m$ Gram matrices, defined by

$$
(K[\mathbf{x}])_{i j}=K\left(x_{i}, x_{j}\right), \quad\left(\hat{K}_{D}[\mathbf{x}]\right)_{i j}=\hat{K}_{D}\left(x_{i}, x_{j}\right), \quad(K[\mathbf{y}])_{i j}=K\left(y_{i}, y_{j}\right), \quad\left(\hat{K}_{D}[\mathbf{y}]\right)_{i j}=\hat{K}_{D}\left(y_{i}, y_{j}\right)
$$

Lemma 6. Let $\mathcal{H}=\ell^{2}$, with $\mathbb{R}^{D}$ identified with a D-dimensional subspace of $\mathcal{H}$ as in Eq. (11). Assume that $\lim _{D \rightarrow \infty} \hat{K}_{D}(x, y)=$ $K(x, y)$ for all pairs $(x, y) \in \mathcal{X} \times \mathcal{X}$. Then

$$
\begin{equation*}
\lim _{D \rightarrow \infty}\left\|C_{\hat{\Phi}_{D}(\mathbf{x})}-C_{\Phi(\mathbf{x})}\right\|_{\mathrm{HS}(\mathcal{H})}=0 \tag{12}
\end{equation*}
$$

Proof. Let $A=\frac{1}{\sqrt{m}} \Phi(\mathbf{x}) J_{m}: \mathbb{R}^{m} \rightarrow \mathcal{H}$, then

$$
A A^{T}=C_{\Phi(\mathbf{x})}, \quad A^{T} A=\frac{1}{m} J_{m} K[\mathbf{x}] J_{m}
$$

By Lemma 5, the nonzero eigenvalues of $C_{\Phi(\mathbf{x})}=A A^{T}$ are the same as those of $\frac{1}{m} J_{m} K[\mathbf{x}] J_{m}=A^{T} A$. Similarly, the nonzero eigenvalues of $C_{\hat{\Phi}_{D}(\mathbf{x})}$ are the same as those of $\frac{1}{m} J_{m} \hat{K}_{D}[\mathbf{x}] J_{m}$. This also implies that both $C_{\Phi(\mathbf{x})}$ and $C_{\hat{\Phi}_{D}(\mathbf{x})}$ have rank at most $m-1$, since $\operatorname{rank}\left(J_{m}\right)=m-1$.

Since $\lim _{D \rightarrow \infty} \hat{K}_{D}\left(x_{i}, x_{j}\right)=K\left(x_{i}, x_{j}\right)$ for all pairs $\left(x_{i}, x_{j}\right), 1 \leq i, j \leq m$, we have, as $m \times m$ matrices,

$$
\lim _{D \rightarrow \infty}\left\|J_{m} \hat{K}_{D}[\mathbf{x}] J_{m}-J_{m} K[\mathbf{x}] J_{m}\right\|_{F}=0
$$

Since $J_{m} \hat{K}_{D}[\mathbf{x}] J_{m}$ and $J_{m} K[\mathbf{x}] J_{m}$ are finite matrices, convergence in the $\left\|\|_{F}\right.$ norm is equivalent to convergence in the operator || || norm. Thus we have

$$
\begin{array}{r}
\lim _{D \rightarrow \infty}\left\|J_{m} \hat{K}_{D}[\mathbf{x}] J_{m}-J_{m} K[\mathbf{x}] J_{m}\right\|=0 \Longleftrightarrow \lim _{D \rightarrow \infty} \lambda_{\max }\left(J_{m} \hat{K}_{D}[\mathbf{x}] J_{m}\right)=\lambda_{\max }\left(J_{m} K[\mathbf{x}] J_{m}\right) \\
\Longleftrightarrow \lim _{D \rightarrow \infty} \lambda_{\max }\left(C_{\hat{\Phi}_{D}(\mathbf{x})}\right)=\lambda_{\max }\left(C_{\Phi(\mathbf{x})}\right) \Longleftrightarrow \lim _{D \rightarrow \infty}\left\|C_{\hat{\Phi}_{D}(\mathbf{x})}-C_{\Phi(\mathbf{x})}\right\|=0 \\
\Longleftrightarrow \lim _{D \rightarrow \infty}\left\|C_{\hat{\Phi}_{D}(\mathbf{x})}-C_{\Phi(\mathbf{x})}\right\|_{\mathrm{HS}(\mathcal{H})}=0
\end{array}
$$

by Lemma 1 , since both $C_{\hat{\Phi}_{D}(\mathbf{x})}, C_{\Phi(\mathbf{x})}$ have rank at most $m-1$. This completes the proof of the lemma.

Proof of Theorem 1. Consider the expansion

$$
\begin{align*}
& \left\|\log \left(C_{\hat{\Phi}_{D}(\mathbf{x})}+\gamma I_{D}\right)-\log \left(C_{\hat{\Phi}_{D}(\mathbf{y})}+\mu I_{D}\right)\right\|_{F}^{2}=\left\|\log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{x})}}{\gamma}+I_{D}\right)-\log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{y})}}{\mu}+I_{D}\right)+(\log \gamma-\log \mu) I_{D}\right\|_{F}^{2} \\
& =\left\|\log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{x})}}{\gamma}+I_{D}\right)-\log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{y})}}{\mu}+I_{D}\right)\right\|_{F}^{2}+2(\log \gamma-\log \mu) \operatorname{tr}\left(\log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{x})}}{\gamma}+I_{D}\right)-\log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{y})}}{\mu}+I_{D}\right)\right) \\
& +(\log \gamma-\log \mu)^{2} D \tag{13}
\end{align*}
$$

With $\mathbb{R}^{D}$ identified as a subspace of $\mathcal{H}=\ell^{2}$, we have by Lemma 6 (with the scaling factors $\gamma, \mu$ ), that

$$
\lim _{D \rightarrow \infty}\left\|\frac{C_{\hat{\Phi}_{D}(\mathbf{x})}}{\gamma}-\frac{C_{\Phi(\mathbf{x})}}{\gamma}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=0, \quad \lim _{D \rightarrow \infty}\left\|\frac{C_{\hat{\Phi}_{D}(\mathbf{y})}}{\mu}-\frac{C_{\Phi(\mathbf{y})}}{\mu}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=0 .
$$

By Lemma 3, we have

$$
\begin{aligned}
& \lim _{D \rightarrow \infty} \operatorname{tr}\left(\log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{x})}}{\gamma}+I_{D}\right)\right)=\operatorname{tr}\left(\log \left(\frac{C_{\Phi(\mathbf{x})}}{\gamma}+I_{\mathcal{H}}\right)\right)=\operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} K[\mathbf{x}] J_{m}+I_{m}\right)\right] \\
& \lim _{D \rightarrow \infty} \operatorname{tr}\left(\log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{y})}}{\mu}+I_{D}\right)\right)=\operatorname{tr}\left(\log \left(\frac{C_{\Phi(\mathbf{y})}}{\mu}+I_{\mathcal{H}}\right)\right)=\operatorname{tr}\left[\log \left(\frac{1}{\mu m} J_{m} K[\mathbf{y}] J_{m}+I_{m}\right)\right]
\end{aligned}
$$

Since these two quantities are both finite, for $\gamma \neq \mu$, as $D \rightarrow \infty$, clearly the right hand side of Eq. (13) goes to infinity. This gives us the desired limit.
Proof of Theorem 2. Without loss of generality, we identify $\mathcal{H}$ with $\ell^{2}$ as above and identify $\mathbb{R}^{D}$ with a $D$-dimensional subspace of $\ell^{2}$ as in Eq. (11). When $\gamma=\mu$, we have

$$
\begin{aligned}
& \left\|\log \left(C_{\Phi(\mathbf{x})}+\gamma I_{\mathcal{H}}\right)-\log \left(C_{\Phi(\mathbf{y})}+\gamma I_{\mathcal{H}}\right)\right\|_{\mathrm{eHS}}^{2}=\left\|\log \left(\frac{C_{\Phi(\mathbf{x})}}{\gamma}+I_{\mathcal{H}}\right)-\log \left(\frac{C_{\Phi(\mathbf{y})}}{\gamma}+I_{\mathcal{H}}\right)\right\|_{\mathrm{HS}}^{2} \\
& =\left\|\log \left(\frac{C_{\Phi(\mathbf{x})}}{\gamma}+I_{\mathcal{H}}\right)\right\|_{\mathrm{HS}}^{2}+\left\|\log \left(\frac{C_{\Phi(\mathbf{y})}}{\gamma}+I_{\mathcal{H}}\right)\right\|_{\mathrm{HS}}^{2}-2 \operatorname{tr}\left[\log \left(\frac{C_{\Phi(\mathbf{x})}}{\gamma}+I_{\mathcal{H}}\right) \log \left(\frac{C_{\Phi(\mathbf{y})}}{\gamma}+I_{\mathcal{H}}\right)\right] .
\end{aligned}
$$

It follows from Lemma 5 that the first term is

$$
\begin{aligned}
\left\|\log \left(\frac{1}{\gamma} C_{\Phi(\mathbf{x})}+I_{\mathcal{H}}\right)\right\|_{\mathrm{HS}}^{2} & =\left\|\log \left(\frac{1}{\gamma m} \Phi(\mathbf{x}) J_{m}^{2} \Phi(\mathbf{x})^{T}+I_{\mathcal{H}}\right)\right\|_{\mathrm{HS}}^{2}=\left\|\log \left(\frac{1}{\gamma m} J_{m} \Phi(\mathbf{x})^{T} \Phi(\mathbf{x}) J_{m}+I_{m}\right)\right\|_{\mathrm{HS}}^{2} \\
& =\left\|\log \left(\frac{1}{\gamma m} J_{m} K[\mathbf{x}] J_{m}+I_{m}\right)\right\|_{\mathrm{HS}}^{2}=\operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} K[\mathbf{x}] J_{m}+I_{m}\right)\right]^{2}
\end{aligned}
$$

Similarly, the second term is

$$
\left\|\log \left(\frac{1}{\gamma} C_{\Phi(\mathbf{y})}+I_{\mathcal{H}}\right)\right\|_{\mathrm{HS}}^{2}=\operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} K[\mathbf{y}] J_{m}+I_{m}\right)\right]^{2}
$$

Thus we have

$$
\begin{align*}
\left\|\log \left(C_{\Phi(\mathbf{x})}+\gamma I_{\mathcal{H}}\right)-\log \left(C_{\Phi(\mathbf{y})}+\gamma I_{\mathcal{H}}\right)\right\|_{\text {eHS }}^{2} & =\operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} K[\mathbf{x}] J_{m}+I_{m}\right)\right]^{2}+\operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} K[\mathbf{y}] J_{m}+I_{m}\right)\right]^{2} \\
& -2 \operatorname{tr}\left[\log \left(\frac{C_{\Phi(\mathbf{x})}}{\gamma}+I_{\mathcal{H}}\right) \log \left(\frac{C_{\Phi(\mathbf{y})}}{\gamma}+I_{\mathcal{H}}\right)\right] \tag{14}
\end{align*}
$$

Similarly,

$$
\begin{align*}
\left\|\log \left(C_{\hat{\Phi}_{D}(\mathbf{x})}+\gamma I_{D}\right)-\log \left(C_{\hat{\Phi}_{D}(\mathbf{y})}+\gamma I_{D}\right)\right\|_{F}^{2} & =\operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} \hat{K}_{D}[\mathbf{x}] J_{m}+I_{m}\right)\right]^{2}+\operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} \hat{K}_{D}[\mathbf{y}] J_{m}+I_{m}\right)\right]^{2} \\
& -2 \operatorname{tr}\left[\log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{x})}}{\gamma}+I_{D}\right) \log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{y})}}{\gamma}+I_{D}\right)\right] . \tag{15}
\end{align*}
$$

With $\mathbb{R}^{D}$ identified as a subspace of $\mathcal{H}=\ell^{2}$, we have by Lemma 6 (with the scaling factor $\gamma$ ), that

$$
\lim _{D \rightarrow \infty}\left\|\frac{C_{\hat{\Phi}_{D}(\mathbf{x})}}{\gamma}-\frac{C_{\Phi(\mathbf{x})}}{\gamma}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=0, \quad \lim _{D \rightarrow \infty}\left\|\frac{C_{\hat{\Phi}_{D}(\mathbf{y})}}{\gamma}-\frac{C_{\Phi(\mathbf{y})}}{\gamma}\right\|_{\mathrm{HS}(\mathcal{H})}^{2}=0,
$$

with the operators $C_{\hat{\Phi}_{D}(\mathbf{x})}, C_{\Phi(\mathbf{x})}, C_{\hat{\Phi}_{D}(\mathbf{y})}, C_{\Phi(\mathbf{y})}$ all have rank at most $m-1$. It thus follows from Lemma 4 that

$$
\begin{equation*}
\lim _{D \rightarrow \infty} \operatorname{tr}\left[\log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{x})}}{\gamma}+I_{D}\right) \log \left(\frac{C_{\hat{\Phi}_{D}(\mathbf{y})}}{\gamma}+I_{D}\right)\right]=\operatorname{tr}\left[\log \left(\frac{C_{\Phi(\mathbf{x})}}{\gamma}+I_{\mathcal{H}}\right) \log \left(\frac{C_{\Phi(\mathbf{y})}}{\gamma}+I_{\mathcal{H}}\right)\right] . \tag{16}
\end{equation*}
$$

Similarly, since $\lim _{D \rightarrow \infty} \hat{K}_{D}\left(x_{i}, x_{j}\right)=K\left(x_{i}, x_{j}\right)$ for all pairs $\left(x_{i}, x_{j}\right)$ and $\lim _{D \rightarrow \infty} \hat{K}_{D}\left(y_{i}, y_{j}\right)=K\left(y_{i}, y_{j}\right)$ for all pairs $\left(y_{i}, y_{j}\right), 1 \leq i, j \leq m$, we have, as $m \times m$ matrices,

$$
\lim _{D \rightarrow \infty}\left\|J_{m} \hat{K}_{D}[\mathbf{x}] J_{m}-J_{m} K[\mathbf{x}] J_{m}\right\|_{F}=0, \quad \lim _{D \rightarrow \infty}\left\|J_{m} \hat{K}_{D}[\mathbf{y}] J_{m}-J_{m} K[\mathbf{y}] J_{m}\right\|_{F}=0
$$

It also follows from Lemma 4 that

$$
\begin{align*}
& \lim _{D \rightarrow \infty} \operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} \hat{K}_{D}[\mathbf{x}] J_{m}+I_{m}\right)\right]^{2}=\operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} K[\mathbf{x}] J_{m}+I_{m}\right)\right]^{2}  \tag{17}\\
& \lim _{D \rightarrow \infty} \operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} \hat{K}_{D}[\mathbf{y}] J_{m}+I_{m}\right)\right]^{2}=\operatorname{tr}\left[\log \left(\frac{1}{\gamma m} J_{m} K[\mathbf{y}] J_{m}+I_{m}\right)\right]^{2} \tag{18}
\end{align*}
$$

Combining the expressions in Eqs. (14), (15), (16), (17), (18), we obtain

$$
\lim _{D \rightarrow \infty}\left\|\log \left(C_{\hat{\Phi}_{D}(\mathbf{x})}+\gamma I_{D}\right)-\log \left(C_{\hat{\Phi}_{D}(\mathbf{y})}+\gamma I_{D}\right)\right\|_{F}^{2}=\left\|\log \left(C_{\Phi(\mathbf{x})}+\gamma I_{\mathcal{H}}\right)-\log \left(C_{\Phi(\mathbf{y})}+\gamma I_{\mathcal{H}}\right)\right\|_{\text {eHS }}^{2}
$$

This completes the proof of the Theorem.

## 2. Further information on the Hilbert-Schmidt distance between covariance operators

For completeness, in this section we provide the mathematical expression for the Hilbert-Schmidt distance between two RKHS covariance operators. This was used for carrying out the corresponding experiments on the Fish dataset in the main paper. Let $K$ be a positive definite kernel on an arbitrary non-empty set $\mathcal{X}$ and $\mathcal{H}_{K}$ be its corresponding RKHS. Let $C_{\Phi(\mathbf{x})}$ and $C_{\Phi(\mathbf{y})}$ be the covariance operators corresponding to two $n \times m$ data matrices $\mathbf{x}$ and $\mathbf{y}$, respectively, sampled from $\mathcal{X}$. Following [2], let $K[\mathbf{x}], K[\mathbf{y}]$, and $K[\mathbf{x}, \mathbf{y}]$ denote the $m \times m$ Gram matrices defined by

$$
(K[\mathbf{x}])_{i j}=K\left(x_{i}, x_{j}\right), \quad(K[\mathbf{y}])_{i j}=K\left(y_{i}, y_{j}\right), \quad(K[\mathbf{x}, \mathbf{y}])_{i j}=K\left(x_{i}, y_{j}\right), \quad 1 \leq i, j \leq m
$$

Then the Gram matrices and the covariance operators are related by

$$
\Phi(\mathbf{x})^{T} \Phi(\mathbf{x})=K[\mathbf{x}], \quad \Phi(\mathbf{y})^{T} \Phi(\mathbf{y})=K[\mathbf{y}], \quad \Phi(\mathbf{x})^{T} \Phi(\mathbf{y})=K[\mathbf{x}, \mathbf{y}] .
$$

Here $\Phi(\mathbf{x})^{T}$ denotes the transpose of $\Phi(\mathbf{x})$ in the case $\operatorname{dim}\left(\mathcal{H}_{K}\right)<\infty$ and the adjoint operator of $\Phi(\mathbf{x})$ in the case $\operatorname{dim}\left(\mathcal{H}_{K}\right)=\infty$.
Lemma 7. The Hilbert-Schmidt distances between two RKHS covariance operators $C_{\Phi(\mathbf{x})}$ and $C_{\Phi(\mathbf{y})}$ is given by

$$
\begin{equation*}
\left\|C_{\Phi(\mathbf{x})}-C_{\Phi(\mathbf{y})}\right\|_{\mathrm{HS}}^{2}=\frac{1}{m^{2}}\left\langle J_{m} K[\mathbf{x}], K[\mathbf{x}] J_{m}\right\rangle_{F}-\frac{2}{m^{2}}\left\langle J_{m} K[\mathbf{x}, \mathbf{y}], K[\mathbf{x}, \mathbf{y}] J_{m}\right\rangle_{F}+\frac{1}{m^{2}}\left\langle J_{m} K[\mathbf{y}], K[\mathbf{y}] J_{m}\right\rangle_{F} \tag{19}
\end{equation*}
$$

Proof of Lemma 7. By definition of the Hilbert-Schmidt norm and property of the trace operation, we have

$$
\begin{aligned}
& \left\|C_{\Phi(\mathbf{x})}-C_{\Phi(\mathbf{y})}\right\|_{\mathrm{HS}}^{2}=\left\|\frac{1}{m} \Phi(\mathbf{x}) J_{m} \Phi(\mathbf{x})^{T}-\frac{1}{m} \Phi(\mathbf{y}) J_{m} \Phi(\mathbf{y})^{T}\right\|_{\mathrm{HS}}^{2} \\
& =\frac{1}{m^{2}}\left\|\Phi(\mathbf{x}) J_{m} \Phi(\mathbf{x})^{T}\right\|_{\mathrm{HS}}^{2}-\frac{2}{m^{2}}\left\langle\Phi(\mathbf{x}) J_{m} \Phi(\mathbf{x})^{T}, \Phi(\mathbf{y}) J_{m} \Phi(\mathbf{y})^{T}\right\rangle_{\mathrm{HS}}+\frac{1}{m^{2}}\left\|\Phi(\mathbf{y}) J_{m} \Phi(\mathbf{y})^{T}\right\|_{\mathrm{HS}}^{2} \\
& =\frac{1}{m^{2}} \operatorname{tr}\left[\Phi(\mathbf{x}) J_{m} \Phi(\mathbf{x})^{T} \Phi(\mathbf{x}) J_{m} \Phi(\mathbf{x})^{T}\right]-\frac{2}{m^{2}} \operatorname{tr}\left[\Phi(\mathbf{x}) J_{m} \Phi(\mathbf{x})^{T} \Phi(\mathbf{y}) J_{m} \Phi(\mathbf{y})^{T}\right]+\frac{1}{m^{2}} \operatorname{tr}\left[\Phi(\mathbf{y}) J_{m} \Phi(\mathbf{y})^{T} \Phi(\mathbf{y}) J_{m} \Phi(\mathbf{y})^{T}\right] \\
& =\frac{1}{m^{2}} \operatorname{tr}\left[\left(K[\mathbf{x}] J_{m}\right)^{2}-2 K[\mathbf{y}, \mathbf{x}] J_{m} K[\mathbf{x}, \mathbf{y}] J_{m}+\left(K[\mathbf{y}] J_{m}\right)^{2}\right] \\
& =\frac{1}{m^{2}}\left[\left\langle J_{m} K[\mathbf{x}], K[\mathbf{x}] J_{m}\right\rangle_{F}-2\left\langle J_{m} K[\mathbf{x}, \mathbf{y}], K[\mathbf{x}, \mathbf{y}] J_{m}\right\rangle_{F}+\left\langle J_{m} K[\mathbf{y}], K[\mathbf{y}] J_{m}\right\rangle_{F}\right] .
\end{aligned}
$$

This completes the proof of the lemma.

## 3. Further information on the Quasi-random Fourier features

Consider again the expression of the kernel $K(x, y)=k(x-y)$ by Bochner's theorem

$$
\begin{align*}
k(x-y) & =\int_{\mathbb{R}^{n}} e^{-i\langle\omega, x-y\rangle} d \rho(\omega)  \tag{20}\\
& =\int_{\mathbb{R}^{n}} \rho(\omega) \phi_{\omega}(x) \overline{\phi_{\omega}(y)} d \omega, \text { where } \phi_{\omega}(x)=e^{-i\langle\omega, x\rangle}
\end{align*}
$$

The Random Fourier feature maps arise from the Monte-Carlo approximation of the integral in Eq. (20), using a random set of points $\omega_{j}$ 's sampled according to the distribution $\rho$. In this section, we describe in more detail the Quasi-random Fourier features approach proposed recently by [4]. This approach is based on the methodology of Quasi-Monte Carlo integration [1], in which the $\omega_{j}$ 's are deterministic points arising from a low-discrepancy sequence in $[0,1]^{n}$ (see below for more details).

Assume that the distribution $\rho$ in Eq. (20) has the product form $\rho(\omega)=\prod_{j=1}^{n} \rho_{j}\left(\omega_{j}\right)$. Assume that each component cumulative distribution function $\psi_{j}\left(x_{j}\right)=\int_{-\infty}^{x_{j}} \rho_{j}\left(z_{j}\right) d z_{j}$ is strictly increasing, so that the inverse functions $\psi_{j}^{-1}:[0,1] \rightarrow$ $\mathbb{R}$ are all well-defined. Let $\psi: \mathbb{R}^{n} \rightarrow[0,1]^{n}$ be defined by $\psi(x)=\psi\left(x_{1}, \ldots, x_{n}\right)=\left(\psi_{1}\left(x_{1}\right), \ldots, \psi_{n}\left(x_{n}\right)\right)$. Then its inverse function $\psi^{-1}:[0,1]^{n} \rightarrow \mathbb{R}^{n}$ is well-defined and is given component-wise by $\psi^{-1}(z)=\psi^{-1}\left(z_{1}, \ldots, z_{n}\right)=$ $\left(\psi_{1}^{-1}\left(z_{1}\right), \ldots, \psi_{n}^{-1}\left(z_{n}\right)\right)$.

With the change of variable $\omega=\psi^{-1}(t)$, the integral in Eq. (20) becomes

$$
\begin{equation*}
\int_{\mathbb{R}^{n}} e^{-i\langle\omega, x-y\rangle} \rho(\omega) d \omega=\int_{[0,1]^{n}} e^{-i\left\langle\psi^{-1}(t), x-y\right\rangle} d t \tag{21}
\end{equation*}
$$

Instead of approximating the left hand side of Eq. (21) using a random set of points $\left\{\omega_{j}\right\}_{j=1}^{D}$ in $\mathbb{R}^{n}$ sampled according to $\rho$, in the Quasi-Monte Carlo approach, one approximates the right hand side using a deterministic, low-discrepancy sequence of points $\left\{t_{j}\right\}_{j=1}^{D}$ in $[0,1]^{n}$. This sequence gives rise to a deterministic sequence

$$
\begin{equation*}
\omega_{j}=\psi^{-1}\left(t_{j}\right), \quad 1 \leq j \leq D \tag{22}
\end{equation*}
$$

from which we construct the Fourier feature map as described by Eqs. (23), (24), and (25),

$$
\begin{align*}
\cos \left(W^{T} x\right) & =\left(\cos \left(\left\langle\omega_{1}, x\right\rangle\right), \ldots, \cos \left(\left\langle\omega_{D}, x\right\rangle\right)\right)^{T} \in \mathbb{R}^{D}  \tag{23}\\
\sin \left(W^{T} x\right) & =\left(\sin \left(\left\langle\omega_{1}, x\right\rangle\right), \ldots, \sin \left(\left\langle\omega_{D}, x\right\rangle\right)\right)^{T} \in \mathbb{R}^{D} .  \tag{24}\\
\hat{\Phi}_{D}(x) & =\frac{1}{\sqrt{D}}\left(\cos \left(W^{T} x\right) ; \sin \left(W^{T} x\right)\right) \in \mathbb{R}^{2 D} \tag{25}
\end{align*}
$$

just as in the case of random Fourier features. In our experiments, $\left\{t_{j}\right\}_{j=1}^{D}$ is a Halton sequence, whose implementation is readily available in MATLAB ${ }^{1}$.

### 3.1. Low-discrepancy sequences

In this section, we briefly review the concept of low-discrepancy sequences in Quasi-Monte Carlo methods. For a comprehensive treatment, we refer to [3]. Let $n \in \mathbb{N}$ be fixed. Let $I^{n}=[0,1)^{n}$ and denote its closure by $\bar{I}^{n}=[0,1]^{n}$. For an integrable function $f$ in $\bar{I}^{n}$, we consider the approximation

$$
\begin{equation*}
\int_{\bar{I}^{n}} f(u) d u \approx \frac{1}{N} \sum_{j=1}^{N} f\left(x_{j}\right) \tag{26}
\end{equation*}
$$

using a deterministic set of points $P=\left(x_{1}, \ldots, x_{N}\right)$, which are part of an infinite sequence $\left(x_{j}\right)_{j \in \mathbb{N}}$ in $\bar{I}^{n}$, such that the integration error satisfies

$$
\begin{equation*}
\lim _{N \rightarrow \infty}\left|\frac{1}{N} \sum_{j=1}^{N} f\left(x_{j}\right)-\int_{\bar{I}^{n}} f(u) d u\right|=0 \tag{27}
\end{equation*}
$$

[^0]This convergence can be measured via the concept of discrepancy as follows. Let $N$ be fixed. For an arbitrary set $B \subset \bar{I}^{n}$, define the counting function

$$
\begin{equation*}
A(B ; P)=\sum_{j=1}^{N} \chi_{B}\left(x_{j}\right) \tag{28}
\end{equation*}
$$

where $\chi_{B}$ denotes the characteristic function for $B$. Thus $A(B ; P)$ denotes the number of points in $P$ that lie in the set $B$.
Let $\mathcal{B}$ be a non-empty family of Lebesgue-measurable subsets of $\bar{I}^{n}$. The discrepancy of the set $P$ with respect to $\mathcal{B}$ is then defined by

$$
\begin{equation*}
D_{N}(\mathcal{B} ; P)=\sup _{B \in \mathcal{B}}\left|\frac{A(B ; P)}{N}-\operatorname{vol}(B)\right|, \tag{29}
\end{equation*}
$$

with $\operatorname{vol}(B)$ denoting the volume of $B$ with respect to the Lebesgue measure.
The star discrepancy $D_{N}^{*}(P)$ is defined by

$$
\begin{equation*}
D_{N}^{*}(P)=D_{N}\left(\mathcal{J}^{*} ; P\right) \tag{30}
\end{equation*}
$$

where $\mathcal{J}^{*}$ denotes the family of all subintervals of $I^{n}$ of the form $\prod_{j=1}^{n}\left[0, x_{j}\right)$. The star discrepancy and the integration error are related via the Koksma- Hlawka inequality, as follows. Define

$$
\begin{equation*}
V(f)=\sum_{k=1}^{n} \sum_{1 \leq i_{1} \leq \cdots \leq i_{k} \leq n} \int_{0}^{1} \cdots \int_{0}^{1}\left|\frac{\partial^{k} f}{\partial u_{i_{1}} \ldots \partial u_{i_{k}}}\right| d u_{i_{1}} \ldots d u_{i_{k}} \tag{31}
\end{equation*}
$$

which is called the variation of $f$ on $\bar{I}^{n}$ in the sense of Hardy-Krause.
Theorem 3 (Koksma-Hlawka inequality). If $f$ has bounded variation $V(f)$ on $\bar{I}^{n}$ in the sense of Hardy-Krause, then for any set $\left(x_{1}, \ldots, x_{N}\right)$ in $I^{n}$,

$$
\begin{equation*}
\left|\frac{1}{N} \sum_{j=1}^{N} f\left(x_{j}\right)-\int_{\bar{I}^{n}} f(u) d u\right| \leq V(f) D_{N}^{*}\left(x_{1}, \ldots, x_{N}\right) \tag{32}
\end{equation*}
$$

By Theorem 3, to achieve a small integration error, we need a sequence $\left(x_{j}\right)_{j \in \mathbb{N}}$ with low discrepancy $D_{N}^{*}\left(x_{1}, \ldots, x_{N}\right) \rightarrow$ 0 as $N \rightarrow \infty$. Some examples of low-discrepancy sequences are Halton and Sobol' sequences (we refer to [3, 1] for the detailed constructions of these and other sequences). The Halton sequence in particular satisfies $D_{N}^{*}\left(x_{1}, \ldots, x_{N}\right)=$ $C(n) \frac{(\log N)^{n}}{N}$ for $N \geq 2$.

### 3.2. The Gaussian case

In this section, we give the explicit expression for the functions $\psi$ and $\psi^{-1}$, as defined above, in the case of the Gaussian kernel. It suffices for us to consider the one-dimensional setting here, since the multivariate case is defined componentwise using the one-dimensional case. For $K(x, y)=e^{-\frac{(x-y)^{2}}{\sigma^{2}}}$, we have $\rho(z)=\frac{\sigma}{2 \sqrt{\pi}} e^{-\frac{\sigma^{2} z^{2}}{4}}$. Recall the Gaussian error function erf defined by $\operatorname{erf}(x)=\frac{2}{\sqrt{\pi}} \int_{0}^{x} e^{-z^{2}} d z$ and the complementary Gaussian error function erfc defined by $\operatorname{erfc}(x)=\frac{2}{\sqrt{\pi}} \int_{x}^{\infty} e^{-z^{2}} d z=1-\operatorname{erf}(x)$. By definition, the cumulative distribution function $\psi$ for $\rho$ is given by

$$
\psi(x)=\int_{-\infty}^{x} \rho(z) d z=1-\int_{x}^{\infty} \rho(z) d z=1-\frac{\sigma}{2 \sqrt{\pi}} \int_{x}^{\infty} e^{-\frac{\sigma^{2} z^{2}}{4}} d z=1-\frac{1}{\sqrt{\pi}} \int_{\frac{x \sigma}{2}}^{\infty} e^{-u^{2}} d u=1-\frac{1}{2} \operatorname{erfc}\left(\frac{x \sigma}{2}\right)
$$

It follows that the inverse function $\psi^{-1}$ is given by

$$
\begin{equation*}
x=\psi^{-1}(t)=\frac{2}{\sigma} \operatorname{erfc}^{-1}(2-2 t)=\frac{2}{\sigma} \operatorname{erf}^{-1}(2 t-1) \tag{33}
\end{equation*}
$$

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[^0]:    ${ }^{1}$ http://www.mathworks.com/help/stats/quasi-random-numbers.html

