Sampling with Stein Discrepancies

Chris. J. Oates

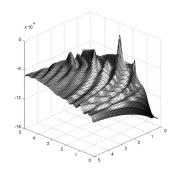
July 2022 @ MCQMC



Motivation: Bayesian Inference and Sampling

The result of integrating expert knowledge with experimental measurement is a *posterior* distribution that is implicitly defined:





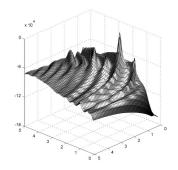
The (lack of) power of computational methodology for Bayesian inference forms a bottleck, imposing a limit of the sophistication of statistical models, or types of information, than can be integrated into a scientific investigation.

This tutorial concerns a new mathematical tool, *Stein discrepancy*, which has the potential to increase the range and scope of statistical analyses that can be performed.

Motivation: Bayesian Inference and Sampling

The result of integrating expert knowledge with experimental measurement is a *posterior* distribution that is implicitly defined:





The (lack of) power of computational methodology for Bayesian inference forms a bottleck, imposing a limit of the sophistication of statistical models, or types of information, than can be integrated into a scientific investigation.

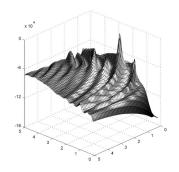
This tutorial concerns a new mathematical tool, *Stein discrepancy*, which has the potential to increase the range and scope of statistical analyses that can be performed.



Motivation: Bayesian Inference and Sampling

The result of integrating expert knowledge with experimental measurement is a *posterior* distribution that is implicitly defined:





The (lack of) power of computational methodology for Bayesian inference forms a bottleck, imposing a limit of the sophistication of statistical models, or types of information, than can be integrated into a scientific investigation.

This tutorial concerns a new mathematical tool, *Stein discrepancy*, which has the potential to increase the range and scope of statistical analyses that can be performed.

Outline of the Tutorial

Measuring Sample Quality with Kernels

Sampling with Kernels

Stein Discrepancy

Case Study: Cardiac Digital Twins
De-Biasing of Markov Chain Monte Carlo

Future Directions, Open Questions and Challenges Scalable Stein Thinning Gradient-Free Kernel Stein Discrepancy Measuring Sample Quality with Kernels

Consider general probability distributions P on general domains \mathcal{X} .

Interested in *sampling*, which in this tutorial means approximation of the form

$$P \approx \sum_{i=1}^n w_i \delta(x_i)$$

for some weights $w_i \in \mathbb{R}$ and some states $x_i \in \mathcal{X}$. [Motivation: Forward UQ.]

The quality of the approximation should be close to optimal (in some sense to be specified) for the number n of states used. [Note: This is not independent sampling from P.]

Markov chain Monte Carlo widely used but "typically" requires $n \approx 10^3$ or 10^4 (to run diagnostics).

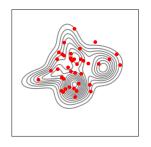
Consider general probability distributions P on general domains \mathcal{X} .

Interested in *sampling*, which in this tutorial means approximation of the form

$$P\approx\sum_{i=1}^n w_i\delta(\mathbf{x}_i)$$

for some weights $w_i \in \mathbb{R}$ and some states $x_i \in \mathcal{X}$.

[Motivation: Forward UQ.]



The quality of the approximation should be close to optimal (in some sense to be specified) for the number n of states used. [Note: This is not independent sampling from P.]

Markov chain Monte Carlo widely used but "typically" requires $n pprox 10^3$ or 10^4 (to run diagnostics).



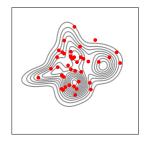
Consider general probability distributions P on general domains \mathcal{X} .

Interested in *sampling*, which in this tutorial means approximation of the form

$$P\approx\sum_{i=1}^n w_i\delta(\mathbf{x}_i)$$

for some weights $w_i \in \mathbb{R}$ and some states $x_i \in \mathcal{X}$.

[Motivation: Forward UQ.]



The quality of the approximation should be close to optimal (in some sense to be specified) for the number n of states used. [Note: This is not independent sampling from P.]

Markov chain Monte Carlo widely used but "typically" requires $n pprox 10^3$ or 10^4 (to run diagnostics).

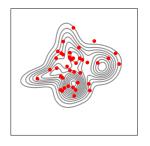
Consider general probability distributions P on general domains \mathcal{X} .

Interested in *sampling*, which in this tutorial means approximation of the form

$$P \approx \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$$

for some weights $w_i \in \mathbb{R}$ and some states $x_i \in \mathcal{X}$.

[Motivation: Forward UQ.]



The quality of the approximation should be close to optimal (in some sense to be specified) for the number n of states used. [Note: This is not independent sampling from P.]

Markov chain Monte Carlo widely used but "typically" requires $npprox 10^3$ or 10^4 (to run diagnostics).

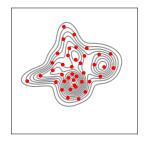
Consider general probability distributions P on general domains \mathcal{X} .

Interested in *sampling*, which in this tutorial means approximation of the form

$$P \approx \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$$

for some weights $w_i \in \mathbb{R}$ and some states $x_i \in \mathcal{X}$.

[Motivation: Forward UQ.]



The quality of the approximation should be close to optimal (in some sense to be specified) for the number n of states used. [Note: This is not independent sampling from P.]

Markov chain Monte Carlo widely used but "typically" requires $n \approx 10^3$ or 10^4 (to run diagnostics).

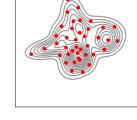


Consider general probability distributions P on general domains \mathcal{X} .

Interested in *sampling*, which in this tutorial means approximation of the form

$$P \approx \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$$

for some weights $w_i \in \mathbb{R}$ and some states $x_i \in \mathcal{X}$. [Motivation: Forward UQ.]



The quality of the approximation should be close to optimal (in some sense to be specified) for the number n of states used. [Note: This is not independent sampling from P.]

Markov chain Monte Carlo widely used but "typically" requires $n \approx 10^3$ or 10^4 (to run diagnostics).



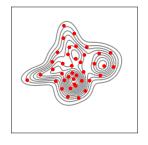
Consider general probability distributions P on general domains \mathcal{X} .

Interested in *sampling*, which in this tutorial means approximation of the form

$$P \approx \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$$

for some weights $w_i \in \mathbb{R}$ and some states $x_i \in \mathcal{X}$.

[Motivation: Forward UQ.]



The quality of the approximation should be close to optimal (in some sense to be specified) for the number n of states used. [Note: This is not independent sampling from P.]

Markov chain Monte Carlo widely used but "typically" requires $n \approx 10^3$ or 10^4 (to run diagnostics).

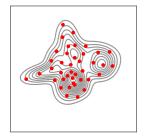


Consider general probability distributions P on general domains \mathcal{X} .

Interested in *sampling*, which in this tutorial means approximation of the form

$$P \approx \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$$

for some weights $w_i \in \mathbb{R}$ and some states $x_i \in \mathcal{X}$. [Motivation: Forward UQ.]



The quality of the approximation should be close to optimal (in some sense to be specified) for the number n of states used. [Note: This is not independent sampling from P.]

Markov chain Monte Carlo widely used but "typically" requires $n \approx 10^3$ or 10^4 (to run diagnostics).



One of the most basic operations that one could hope to perform with a probability distribution is to compute expectations; i.e. to compute integrals of the form

$$\int f dP$$
, or $\int f(x)p(x)dx$

if the probability distribution P admits a probability density function (PDF) p(x).

For implicitly defined P, such integrals do not possess a closed form and numerical integration (also called *cubature*) will be required.

Using a sample, we arrive at a natural approximation

$$P pprox \sum_{i=1}^{n} w_i \delta(x_i) \qquad \Longrightarrow \qquad \int f dP pprox \sum_{i=1}^{n} w_i f(x_i)$$

to the integral. This suggests designing a sample for which the cubature error

$$\int f dP - \sum_{i=1}^n w_i f(x_i)$$

is small. [Problem: Small cubature error could happen by "chance".



One of the most basic operations that one could hope to perform with a probability distribution is to compute expectations; i.e. to compute integrals of the form

$$\int f dP$$
, or $\int f(x)p(x)dx$

if the probability distribution P admits a PDF p(x).

For implicitly defined P, such integrals do not possess a closed form and numerical integration (also called *cubature*) will be required.

Using a sample, we arrive at a natural approximation

$$P pprox \sum_{i=1}^{n} w_i \delta(\mathbf{x}_i) \qquad \Longrightarrow \qquad \int f dP pprox \sum_{i=1}^{n} w_i f(\mathbf{x}_i)$$

to the integral. This suggests designing a sample for which the cubature error

$$\int f dP - \sum_{i=1}^n w_i f(x_i)$$

is small. [Problem: Small cubature error could happen by "chance".



One of the most basic operations that one could hope to perform with a probability distribution is to compute expectations; i.e. to compute integrals of the form

$$\int f dP$$
, or $\int f(x)p(x)dx$

if the probability distribution P admits a PDF p(x).

For implicitly defined P, such integrals do not possess a closed form and numerical integration (also called *cubature*) will be required.

Using a sample, we arrive at a natural approximation

$$P pprox \sum_{i=1}^{n} w_i \delta(\mathbf{x}_i) \qquad \Longrightarrow \qquad \int f dP pprox \sum_{i=1}^{n} w_i f(\mathbf{x}_i)$$

to the integral. This suggests designing a sample for which the cubature error

$$\int f dP - \sum_{i=1}^{n} w_i f(x_i)$$

is small. [Problem: Small cubature error could happen by "chance".



One of the most basic operations that one could hope to perform with a probability distribution is to compute expectations; i.e. to compute integrals of the form

$$\int f dP$$
, or $\int f(x)p(x)dx$

if the probability distribution P admits a PDF p(x).

For implicitly defined P, such integrals do not possess a closed form and numerical integration (also called *cubature*) will be required.

Using a sample, we arrive at a natural approximation

$$P pprox \sum_{i=1}^{n} w_i \delta(\mathbf{x}_i) \qquad \Longrightarrow \qquad \int f dP pprox \sum_{i=1}^{n} w_i f(\mathbf{x}_i)$$

to the integral. This suggests designing a sample for which the cubature error

$$\int f \mathrm{d}P - \sum_{i=1}^n w_i f(\mathbf{x}_i)$$

is small. [Problem: Small cubature error could happen by "chance".]



One of the most basic operations that one could hope to perform with a probability distribution is to compute expectations; i.e. to compute integrals of the form

$$\int f dP$$
, or $\int f(x)p(x)dx$

if the probability distribution P admits a PDF p(x).

For implicitly defined P, such integrals do not possess a closed form and numerical integration (also called *cubature*) will be required.

Using a sample, we arrive at a natural approximation

$$P pprox \sum_{i=1}^{n} w_i \delta(\mathbf{x}_i) \qquad \Longrightarrow \qquad \int f dP pprox \sum_{i=1}^{n} w_i f(\mathbf{x}_i)$$

to the integral. This suggests designing a sample for which the cubature error

$$\int f \mathrm{d}P - \sum_{i=1}^n w_i f(\mathbf{x}_i)$$

is small. [Problem: Small cubature error could happen by "chance".]



The basic idea is as follows: we consider the set S(k) of all functions of the form

$$f(\mathbf{x}) = \sum_{i=1}^m b_i k(\mathbf{x}, \mathbf{y}_i),$$

where k is to be specified, the y_i are fixed states, and $n \in \mathbb{N}$. The function k determines the regularity of the elements in S(k):

Figure: The left panel represents functions built from $k(x,y) = \exp(-\|x-y\|)$, which is non-differentiable, while the right panel corresponds to functions built from $k(x,y) = \exp(-\|x-y\|^2)$, which is smooth.

The basic idea is as follows: we consider the set S(k) of all functions of the form

$$f(\mathbf{x}) = \sum_{i=1}^m b_i k(\mathbf{x}, \mathbf{y}_i),$$

where k is to be specified, the y_i are fixed states, and $n \in \mathbb{N}$. The function k determines the regularity of the elements in S(k):

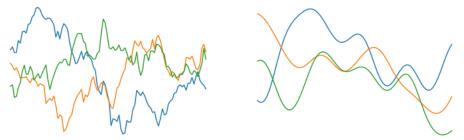


Figure: The left panel represents functions built from $k(x,y) = \exp(-\|x-y\|)$, which is non-differentiable, while the right panel corresponds to functions built from $k(x,y) = \exp(-\|x-y\|^2)$, which is smooth.

S(k) is a vector space (over the reals) of functions when (pointwise) addition and scalar multiplication are defined. In addition to that, we will want to make use of an inner product

$$\langle f,g\rangle_{\mathcal{S}(k)} = \sum_{i=1}^m \sum_{j=1}^n b_i c_j k(\mathbf{y}_i,\mathbf{z}_j), \quad f(\mathbf{x}) = \sum_{i=1}^m b_i k(\mathbf{x},\mathbf{y}_i), \quad g(\mathbf{x}) = \sum_{j=1}^n c_j k(\mathbf{x},\mathbf{z}_j),$$

for which we must require that k is symmetric (i.e. $\langle f, g \rangle_{S(k)} = \langle g, f \rangle_{S(k)}$) and positive definite (i.e. $\langle f, f \rangle_{S(k)} > 0$ for all $f \neq 0$). This inner product is useful because it satisfies a reproducing property, meaning that

$$\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{S}(k)} = f(\mathbf{x}),$$

and suggesting the formal manipulation

$$\frac{1}{n} \sum_{i=1}^{n} f(x_i) - \int f(x) dP(x) = \frac{1}{n} \sum_{i=1}^{n} \langle f, k(\cdot, x_i) \rangle_{S(k)} - \int \langle f, k(\cdot, x) \rangle_{S(k)} dP(x)
= \left\langle f, \underbrace{\frac{1}{n} \sum_{i=1}^{n} k(\cdot, x_i) - \int k(\cdot, x) dP(x)}_{S(k)} \right\rangle_{S(k)},$$

where the cubature error representer $e(\cdot)$ is independent of f. Idea: Seek sample with $\|e\|_{\mathcal{S}(k)}$ small



 $\mathcal{S}(k)$ is a vector space (over the reals) of functions when (pointwise) addition and scalar multiplication are defined. In addition to that, we will want to make use of an inner product

$$\langle f,g\rangle_{\mathcal{S}(k)}=\sum_{i=1}^m\sum_{j=1}^nb_ic_jk(\mathbf{y}_i,\mathbf{z}_j),\quad f(\mathbf{x})=\sum_{i=1}^mb_ik(\mathbf{x},\mathbf{y}_i),\quad g(\mathbf{x})=\sum_{j=1}^nc_jk(\mathbf{x},\mathbf{z}_j),$$

for which we must require that k is symmetric (i.e. $\langle f,g\rangle_{\mathcal{S}(k)}=\langle g,f\rangle_{\mathcal{S}(k)}$) and positive definite (i.e. $\langle f,f\rangle_{\mathcal{S}(k)}>0$ for all $f\neq 0$). This inner product is useful because it satisfies a reproducing property, meaning that

$$\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{S}(k)} = f(\mathbf{x}),$$

and suggesting the formal manipulation

$$\frac{1}{n} \sum_{i=1}^{n} f(x_i) - \int f(x) dP(x) = \frac{1}{n} \sum_{i=1}^{n} \langle f, k(\cdot, x_i) \rangle_{S(k)} - \int \langle f, k(\cdot, x) \rangle_{S(k)} dP(x)
= \left\langle f, \underbrace{\frac{1}{n} \sum_{i=1}^{n} k(\cdot, x_i) - \int k(\cdot, x) dP(x)}_{S(k)} \right\rangle_{S(k)},$$

where the cubature error representer $e(\cdot)$ is independent of f. Idea: Seek sample with $\|e\|_{\mathcal{S}(k)}$ small



S(k) is a vector space (over the reals) of functions when (pointwise) addition and scalar multiplication are defined. In addition to that, we will want to make use of an inner product

$$\langle f,g\rangle_{\mathcal{S}(k)}=\sum_{i=1}^m\sum_{j=1}^nb_ic_jk(y_i,z_j),\quad f(x)=\sum_{i=1}^mb_ik(x,y_i),\quad g(x)=\sum_{j=1}^nc_jk(x,z_j),$$

for which we must require that k is symmetric (i.e. $\langle f,g\rangle_{\mathcal{S}(k)}=\langle g,f\rangle_{\mathcal{S}(k)}$) and positive definite (i.e. $\langle f,f\rangle_{\mathcal{S}(k)}>0$ for all $f\neq 0$). This inner product is useful because it satisfies a reproducing property, meaning that

$$\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{S}(k)} = f(\mathbf{x}),$$

and suggesting the formal manipulation

$$\frac{1}{n} \sum_{i=1}^{n} f(\mathbf{x}_{i}) - \int f(\mathbf{x}) dP(\mathbf{x}) = \frac{1}{n} \sum_{i=1}^{n} \langle f, k(\cdot, \mathbf{x}_{i}) \rangle_{S(k)} - \int \langle f, k(\cdot, \mathbf{x}) \rangle_{S(k)} dP(\mathbf{x})$$

$$= \left\langle f, \underbrace{\frac{1}{n} \sum_{i=1}^{n} k(\cdot, \mathbf{x}_{i}) - \int k(\cdot, \mathbf{x}) dP(\mathbf{x})}_{S(k)} \right\rangle_{S(k)},$$

where the cubature error representer $e(\cdot)$ is independent of f. Idea: Seek sample with $\|e\|_{S(k)}$ small.



S(k) is a vector space (over the reals) of functions when (pointwise) addition and scalar multiplication are defined. In addition to that, we will want to make use of an inner product

$$\langle f,g\rangle_{\mathcal{S}(k)}=\sum_{i=1}^m\sum_{j=1}^nb_ic_jk(\mathbf{y}_i,\mathbf{z}_j),\quad f(\mathbf{x})=\sum_{i=1}^mb_ik(\mathbf{x},\mathbf{y}_i),\quad g(\mathbf{x})=\sum_{j=1}^nc_jk(\mathbf{x},\mathbf{z}_j),$$

for which we must require that k is symmetric (i.e. $\langle f,g\rangle_{\mathcal{S}(k)}=\langle g,f\rangle_{\mathcal{S}(k)}$) and positive definite (i.e. $\langle f,f\rangle_{\mathcal{S}(k)}>0$ for all $f\neq 0$). This inner product is useful because it satisfies a reproducing property, meaning that

$$\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{S}(k)} = f(\mathbf{x}),$$

and suggesting the formal manipulation

$$\frac{1}{n}\sum_{i=1}^{n}f(\mathbf{x}_{i})-\int f(\mathbf{x})\mathrm{d}P(\mathbf{x})=\frac{1}{n}\sum_{i=1}^{n}\langle f,k(\cdot,\mathbf{x}_{i})\rangle_{\mathcal{S}(k)}-\int \langle f,k(\cdot,\mathbf{x})\rangle_{\mathcal{S}(k)}\mathrm{d}P(\mathbf{x})$$

$$=\left\langle f,\frac{1}{n}\sum_{i=1}^{n}k(\cdot,\mathbf{x}_{i})-\int k(\cdot,\mathbf{x})\mathrm{d}P(\mathbf{x})\right\rangle_{\mathcal{S}(k)},$$

where the cubature error representer $e(\cdot)$ is independent of f. Idea: Seek sample with $\|e\|_{S(k)}$ small.



S(k) is a vector space (over the reals) of functions when (pointwise) addition and scalar multiplication are defined. In addition to that, we will want to make use of an inner product

$$\langle f,g\rangle_{\mathcal{S}(k)}=\sum_{i=1}^m\sum_{j=1}^nb_ic_jk(y_i,z_j),\quad f(x)=\sum_{i=1}^mb_ik(x,y_i),\quad g(x)=\sum_{j=1}^nc_jk(x,z_j),$$

for which we must require that k is symmetric (i.e. $\langle f,g\rangle_{\mathcal{S}(k)}=\langle g,f\rangle_{\mathcal{S}(k)}$) and positive definite (i.e. $\langle f,f\rangle_{\mathcal{S}(k)}>0$ for all $f\neq 0$). This inner product is useful because it satisfies a reproducing property, meaning that

$$\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{S}(k)} = f(\mathbf{x}),$$

and suggesting the formal manipulation

$$\frac{1}{n}\sum_{i=1}^{n}f(\mathbf{x}_{i})-\int f(\mathbf{x})\mathrm{d}P(\mathbf{x})=\frac{1}{n}\sum_{i=1}^{n}\langle f,k(\cdot,\mathbf{x}_{i})\rangle_{\mathcal{S}(k)}-\int \langle f,k(\cdot,\mathbf{x})\rangle_{\mathcal{S}(k)}\mathrm{d}P(\mathbf{x})$$

$$=\left\langle f,\underbrace{\frac{1}{n}\sum_{i=1}^{n}k(\cdot,\mathbf{x}_{i})-\int k(\cdot,\mathbf{x})\mathrm{d}P(\mathbf{x})}_{\mathcal{S}(k)}\right\rangle_{\mathcal{S}(k)},$$

where the cubature error representer $e(\cdot)$ is independent of f. Idea: Seek sample with $\|e\|_{S(k)}$ small.



S(k) is a vector space (over the reals) of functions when (pointwise) addition and scalar multiplication are defined. In addition to that, we will want to make use of an inner product

$$\langle f,g\rangle_{\mathcal{S}(k)}=\sum_{i=1}^m\sum_{j=1}^nb_ic_jk(y_i,z_j),\quad f(x)=\sum_{i=1}^mb_ik(x,y_i),\quad g(x)=\sum_{j=1}^nc_jk(x,z_j),$$

for which we must require that k is symmetric (i.e. $\langle f,g\rangle_{\mathcal{S}(k)}=\langle g,f\rangle_{\mathcal{S}(k)}$) and positive definite (i.e. $\langle f,f\rangle_{\mathcal{S}(k)}>0$ for all $f\neq 0$). This inner product is useful because it satisfies a reproducing property, meaning that

$$\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{S}(k)} = f(\mathbf{x}),$$

and suggesting the formal manipulation

$$\frac{1}{n} \sum_{i=1}^{n} f(\mathbf{x}_{i}) - \int f(\mathbf{x}) dP(\mathbf{x}) = \frac{1}{n} \sum_{i=1}^{n} \langle f, k(\cdot, \mathbf{x}_{i}) \rangle_{S(k)} - \int \langle f, k(\cdot, \mathbf{x}) \rangle_{S(k)} dP(\mathbf{x})$$

$$= \left\langle f, \underbrace{\frac{1}{n} \sum_{i=1}^{n} k(\cdot, \mathbf{x}_{i}) - \int k(\cdot, \mathbf{x}) dP(\mathbf{x})}_{S(k)} \right\rangle_{S(k)},$$

where the *cubature error representer* $e(\cdot)$ is indepenent of f. Idea: Seek sample with $\|e\|_{S(k)}$ small.



S(k) is a vector space (over the reals) of functions when (pointwise) addition and scalar multiplication are defined. In addition to that, we will want to make use of an inner product

$$\langle f,g\rangle_{\mathcal{S}(k)}=\sum_{i=1}^m\sum_{j=1}^nb_ic_jk(y_i,z_j),\quad f(x)=\sum_{i=1}^mb_ik(x,y_i),\quad g(x)=\sum_{j=1}^nc_jk(x,z_j),$$

for which we must require that k is symmetric (i.e. $\langle f,g\rangle_{\mathcal{S}(k)}=\langle g,f\rangle_{\mathcal{S}(k)}$) and positive definite (i.e. $\langle f,f\rangle_{\mathcal{S}(k)}>0$ for all $f\neq 0$). This inner product is useful because it satisfies a reproducing property, meaning that

$$\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{S}(k)} = f(\mathbf{x}),$$

and suggesting the formal manipulation

$$\frac{1}{n}\sum_{i=1}^{n}f(\mathbf{x}_{i})-\int f(\mathbf{x})\mathrm{d}P(\mathbf{x})=\frac{1}{n}\sum_{i=1}^{n}\langle f,k(\cdot,\mathbf{x}_{i})\rangle_{\mathcal{S}(k)}-\int \langle f,k(\cdot,\mathbf{x})\rangle_{\mathcal{S}(k)}\mathrm{d}P(\mathbf{x})$$

$$=\left\langle f,\underbrace{\frac{1}{n}\sum_{i=1}^{n}k(\cdot,\mathbf{x}_{i})-\int k(\cdot,\mathbf{x})\mathrm{d}P(\mathbf{x})}_{\mathcal{S}(k)}\right\rangle_{\mathcal{S}(k)},$$

where the cubature error representer $e(\cdot)$ is independent of f. Idea: Seek sample with $\|e\|_{\mathcal{S}(k)}$ small.



S(k) is a vector space (over the reals) of functions when (pointwise) addition and scalar multiplication are defined. In addition to that, we will want to make use of an inner product

$$\langle f, g \rangle_{S(k)} = \sum_{i=1}^{m} \sum_{j=1}^{n} b_i c_j k(y_i, z_j), \quad f(x) = \sum_{i=1}^{m} b_i k(x, y_i), \quad g(x) = \sum_{j=1}^{n} c_j k(x, z_j),$$

for which we must require that k is symmetric (i.e. $\langle f,g\rangle_{\mathcal{S}(k)}=\langle g,f\rangle_{\mathcal{S}(k)}$) and positive definite (i.e. $\langle f,f\rangle_{\mathcal{S}(k)}>0$ for all $f\neq 0$). This inner product is useful because it satisfies a reproducing property, meaning that

$$\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{S}(k)} = f(\mathbf{x}),$$

and suggesting the formal manipulation

$$\frac{1}{n} \sum_{i=1}^{n} f(\mathbf{x}_{i}) - \int f(\mathbf{x}) dP(\mathbf{x}) = \frac{1}{n} \sum_{i=1}^{n} \langle f, k(\cdot, \mathbf{x}_{i}) \rangle_{S(k)} - \int \langle f, k(\cdot, \mathbf{x}) \rangle_{S(k)} dP(\mathbf{x})$$

$$= \left\langle f, \underbrace{\frac{1}{n} \sum_{i=1}^{n} k(\cdot, \mathbf{x}_{i}) - \int k(\cdot, \mathbf{x}) dP(\mathbf{x})}_{S(k)} \right\rangle_{S(k)},$$

where the cubature error representer $e(\cdot)$ is indepenent of f. Idea: Seek sample with $\|e\|_{\mathcal{S}(k)}$ small.



The main mathematical tool that we will exploit is that of a *reproducing kernel*:

Definition 1 (Reproducing kernel Hilbert space)

- 1. $k(\cdot, x) \in \mathcal{H}(k)$ for all $x \in \mathcal{X}$
- 2. $\langle f, k(\cdot, x) \rangle_{\mathcal{H}(k)} = f(x)$ for all $x \in \mathcal{X}$ and all $f \in \mathcal{H}(k)$.
- ▶ Given a symmetric positive definite function k, it can be shown that there exists a unique RKHS $\mathcal{H}(k)$.
- Conversely, each RKHS admits a unique reproducing kernel, and that kernel is symmetric and positive definite.

The main mathematical tool that we will exploit is that of a *reproducing kernel*:

Definition 1 (Reproducing kernel Hilbert space)

- 1. $k(\cdot, x) \in \mathcal{H}(k)$ for all $x \in \mathcal{X}$
- 2. $\langle f, k(\cdot, x) \rangle_{\mathcal{H}(k)} = f(x)$ for all $x \in \mathcal{X}$ and all $f \in \mathcal{H}(k)$.
- ▶ Given a symmetric positive definite function k, it can be shown that there exists a unique RKHS $\mathcal{H}(k)$.
- Conversely, each RKHS admits a unique reproducing kernel, and that kernel is symmetric and positive definite.

The main mathematical tool that we will exploit is that of a reproducing kernel:

Definition 1 (Reproducing kernel Hilbert space)

- 1. $k(\cdot, x) \in \mathcal{H}(k)$ for all $x \in \mathcal{X}$
- 2. $\langle f, k(\cdot, x) \rangle_{\mathcal{H}(k)} = f(x)$ for all $x \in \mathcal{X}$ and all $f \in \mathcal{H}(k)$.
- ▶ Given a symmetric positive definite function k, it can be shown that there exists a unique RKHS $\mathcal{H}(k)$.
- Conversely, each RKHS admits a unique reproducing kernel, and that kernel is symmetric and positive definite.

The main mathematical tool that we will exploit is that of a *reproducing kernel*:

Definition 1 (Reproducing kernel Hilbert space)

- 1. $k(\cdot, x) \in \mathcal{H}(k)$ for all $x \in \mathcal{X}$
- 2. $\langle f, k(\cdot, x) \rangle_{\mathcal{H}(k)} = f(x)$ for all $x \in \mathcal{X}$ and all $f \in \mathcal{H}(k)$.
- ▶ Given a symmetric positive definite function k, it can be shown that there exists a unique RKHS $\mathcal{H}(k)$.
- Conversely, each RKHS admits a unique reproducing kernel, and that kernel is symmetric and positive definite.

The main mathematical tool that we will exploit is that of a reproducing kernel:

Definition 1 (Reproducing kernel Hilbert space)

- 1. $k(\cdot, x) \in \mathcal{H}(k)$ for all $x \in \mathcal{X}$
- 2. $\langle f, k(\cdot, x) \rangle_{\mathcal{H}(k)} = f(x)$ for all $x \in \mathcal{X}$ and all $f \in \mathcal{H}(k)$.
- ▶ Given a symmetric positive definite function k, it can be shown that there exists a unique RKHS $\mathcal{H}(k)$.
- Conversely, each RKHS admits a unique reproducing kernel, and that kernel is symmetric and positive definite.

The main mathematical tool that we will exploit is that of a reproducing kernel:

Definition 1 (Reproducing kernel Hilbert space)

- 1. $k(\cdot, \mathbf{x}) \in \mathcal{H}(k)$ for all $\mathbf{x} \in \mathcal{X}$
- 2. $\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k)} = f(\mathbf{x})$ for all $\mathbf{x} \in \mathcal{X}$ and all $f \in \mathcal{H}(k)$.
- Given a symmetric positive definite function k, it can be shown that there exists a unique RKHS $\mathcal{H}(k)$.
- Conversely, each RKHS admits a unique reproducing kernel, and that kernel is symmetric and positive definite.

The main mathematical tool that we will exploit is that of a reproducing kernel:

Definition 1 (Reproducing kernel Hilbert space)

- 1. $k(\cdot, x) \in \mathcal{H}(k)$ for all $x \in \mathcal{X}$
- 2. $\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k)} = f(\mathbf{x})$ for all $\mathbf{x} \in \mathcal{X}$ and all $f \in \mathcal{H}(k)$.
- Given a symmetric positive definite function k, it can be shown that there exists a unique RKHS $\mathcal{H}(k)$.
- Conversely, each RKHS admits a unique reproducing kernel, and that kernel is symmetric and positive definite.

In general it is difficult to characterise the inner product induced by a reproducing kernel, and hence the elements of the RKHS. However, there are a number of important cases where this can be carried out:

Example 2

The linear span of a finite collection of functions $e_1(x), \ldots, e_p(x)$ can be endowed with the structure of an RKHS with reproducing kernel

$$k(x,y) = \sum_{i=1}^{p} e_i(x)e_i(y).$$

The induced inner product is

$$\langle f, g \rangle_{\mathcal{H}(k)} = b_1 c_1 + \dots + b_p c_p, \qquad f(x) = \sum_{i=1}^n b_i e_i(x), \qquad g(x) = \sum_{i=1}^n c_i e_i(x)$$

Example 3

The kerne

$$k(x, y) = \prod_{i=1}^{d} (1 + \min(1 - x_i, 1 - y_i)), \quad x, y \in [0, 1]^d$$

reproduces a Hilbert space with inner product

$$\langle f,g\rangle_{\mathcal{H}(k)} = \sum_{\mathfrak{u} \subseteq \{1,\ldots,d\}} \int_{[0,1]^d} \frac{\partial^{|\mathfrak{u}|} f}{\partial x_\mathfrak{u}}(x_\mathfrak{u},1) \frac{\partial^{|\mathfrak{u}|} g}{\partial x_\mathfrak{u}}(x_\mathfrak{u},1) \mathrm{d} x_\mathfrak{u}$$



In general it is difficult to characterise the inner product induced by a reproducing kernel, and hence the elements of the RKHS. However, there are a number of important cases where this can be carried out:

Example 2

The linear span of a finite collection of functions $e_1(x), \ldots, e_p(x)$ can be endowed with the structure of an RKHS with reproducing kernel

$$k(\mathbf{x},\mathbf{y})=\sum_{i=1}^p e_i(\mathbf{x})e_i(\mathbf{y}).$$

The induced inner product is

$$\langle f, g \rangle_{\mathcal{H}(k)} = b_1 c_1 + \dots + b_p c_p, \qquad f(x) = \sum_{i=1}^n b_i e_i(x), \qquad g(x) = \sum_{i=1}^n c_i e_i(x)$$

Example 3

The kerne

$$k(x, y) = \prod_{i=1}^{d} (1 + \min(1 - x_i, 1 - y_i)), \quad x, y \in [0, 1]^d$$

reproduces a Hilbert space with inner product

$$\langle f,g\rangle_{\mathcal{H}(k)} = \sum_{\mathfrak{u}\subseteq\{1,\ldots,d\}} \int_{[0,1]^d} \frac{\partial^{|\mathfrak{u}|}f}{\partial x_\mathfrak{u}}(x_\mathfrak{u},1) \frac{\partial^{|\mathfrak{u}|}g}{\partial x_\mathfrak{u}}(x_\mathfrak{u},1) \mathrm{d}x_\mathfrak{u}$$



In general it is difficult to characterise the inner product induced by a reproducing kernel, and hence the elements of the RKHS. However, there are a number of important cases where this can be carried out:

Example 2

The linear span of a finite collection of functions $e_1(x), \ldots, e_p(x)$ can be endowed with the structure of an RKHS with reproducing kernel

$$k(\mathbf{x},\mathbf{y})=\sum_{i=1}^p e_i(\mathbf{x})e_i(\mathbf{y}).$$

The induced inner product is

$$\langle f,g\rangle_{\mathcal{H}(k)}=b_1c_1+\cdots+b_pc_p, \qquad f(x)=\sum_{i=1}^nb_ie_i(x), \qquad g(x)=\sum_{i=1}^nc_ie_i(x).$$

Example

The kerne

$$k(x, y) = \prod_{i=1}^{d} (1 + \min(1 - x_i, 1 - y_i)), \quad x, y \in [0, 1]^d$$

reproduces a Hilbert space with inner product

$$\langle f,g\rangle_{\mathcal{H}(k)} = \sum_{\mathfrak{u}\subseteq\{1,\ldots,d\}} \int_{[0,1]^d} \frac{\partial^{|\mathfrak{u}|}f}{\partial x_\mathfrak{u}}(x_\mathfrak{u},1) \frac{\partial^{|\mathfrak{u}|}g}{\partial x_\mathfrak{u}}(x_\mathfrak{u},1) \mathrm{d}x_\mathfrak{u}$$



In general it is difficult to characterise the inner product induced by a reproducing kernel, and hence the elements of the RKHS. However, there are a number of important cases where this can be carried out:

Example 2

The linear span of a finite collection of functions $e_1(x), \ldots, e_p(x)$ can be endowed with the structure of an RKHS with reproducing kernel

$$k(\mathbf{x},\mathbf{y})=\sum_{i=1}^p e_i(\mathbf{x})e_i(\mathbf{y}).$$

The induced inner product is

$$\langle f,g\rangle_{\mathcal{H}(k)}=b_1c_1+\cdots+b_pc_p, \qquad f(x)=\sum_{i=1}^nb_ie_i(x), \qquad g(x)=\sum_{i=1}^nc_ie_i(x).$$

Example 3

The kernel

$$k(x, y) = \prod_{i=1}^{d} (1 + \min(1 - x_i, 1 - y_i)), \quad x, y \in [0, 1]^d$$

reproduces a Hilbert space with inner product

$$\langle f, g \rangle_{\mathcal{H}(k)} = \sum_{\mathfrak{u} \subset \{1, \dots, d\}} \int_{[0,1]^d} \frac{\partial^{|\mathfrak{u}|} f}{\partial x_{\mathfrak{u}}} (x_{\mathfrak{u}}, 1) \frac{\partial^{|\mathfrak{u}|} g}{\partial x_{\mathfrak{u}}} (x_{\mathfrak{u}}, 1) dx_{\mathfrak{u}}$$



Idea: Work with $\mathcal{H}(k)$ instead of $\mathcal{S}(k)$.

Definition 4 (Kernel mean embedding)

For a kernel k and a probability distribution P, we call $\mu_P = \int k(\cdot, x) dP(x)$ the kernel mean embedding of P in $\mathcal{H}(k)$, whenever it is well-defined.

Example 5

For the kernel

$$k(x,y) = \sum_{i=1}^{\rho} e_i(x)e_i(y),$$

with $e_i \in L^1(P)$, we have kernel mean embedding

$$\mu_P(x) = \sum_{i=1}^p \left(\underbrace{\int e_i(y) \, \mathrm{d}P(y)}_{\leq \infty} \right) e_i(x) \in \mathcal{H}(k).$$

Idea: Work with $\mathcal{H}(k)$ instead of $\mathcal{S}(k)$.

Definition 4 (Kernel mean embedding)

For a kernel k and a probability distribution P, we call $\mu_P = \int k(\cdot, x) dP(x)$ the kernel mean embedding of P in $\mathcal{H}(k)$, whenever it is well-defined.

Example 5

For the kerne

$$k(x,y) = \sum_{i=1}^{p} e_i(x)e_i(y),$$

with $e_i \in L^1(P)$, we have kernel mean embedding

$$\mu_{P}(x) = \sum_{i=1}^{p} \left(\underbrace{\int e_{i}(y) dP(y)}_{\leq \infty} \right) e_{i}(x) \in \mathcal{H}(k)$$



Idea: Work with $\mathcal{H}(k)$ instead of $\mathcal{S}(k)$.

Definition 4 (Kernel mean embedding)

For a kernel k and a probability distribution P, we call $\mu_P = \int k(\cdot, x) dP(x)$ the kernel mean embedding of P in $\mathcal{H}(k)$, whenever it is well-defined.

Example 5

For the kernel

$$k(\mathbf{x},\mathbf{y})=\sum_{i=1}^p e_i(\mathbf{x})e_i(\mathbf{y}),$$

with $e_i \in L^1(P)$, we have kernel mean embedding

$$\mu_P(\mathbf{x}) = \sum_{i=1}^p \left(\underbrace{\int e_i(\mathbf{y}) dP(\mathbf{y})}_{\leq \infty}\right) e_i(\mathbf{x}) \in \mathcal{H}(k).$$



Idea: Work with $\mathcal{H}(k)$ instead of $\mathcal{S}(k)$.

Definition 4 (Kernel mean embedding)

For a kernel k and a probability distribution P, we call $\mu_P = \int k(\cdot, x) dP(x)$ the kernel mean embedding of P in $\mathcal{H}(k)$, whenever it is well-defined.

Example 5

For the kernel

$$k(\mathbf{x},\mathbf{y})=\sum_{i=1}^p e_i(\mathbf{x})e_i(\mathbf{y}),$$

with $e_i \in L^1(P)$, we have kernel mean embedding

$$\mu_P(\mathbf{x}) = \sum_{i=1}^p \left(\underbrace{\int e_i(\mathbf{y}) dP(\mathbf{y})}_{\leq \infty}\right) e_i(\mathbf{x}) \in \mathcal{H}(k).$$

Lemma 6

If $\int \sqrt{k(\mathbf{x},\mathbf{x})} \mathrm{d}P(\mathbf{x}) < \infty$ then $\mu_P(\mathbf{x}) \in \mathcal{H}(k)$.

Proof

Consider the linear operator $Lf = \int f(x) dP(x)$ acting on $f \in \mathcal{H}(k)$. Claim that L is a bounded linear operator from $\mathcal{H}(k)$ to \mathbb{R} . Indeed,

$$|Lf| = \left| \int f(x) dP(x) \right| \le \int |f(x)| dP(x) \tag{1}$$

$$= \int |\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k)} | \mathrm{d}P(\mathbf{x})$$
 (2)

$$\leq \int \|f\|_{\mathcal{H}(k)} \|k(\cdot, \mathbf{x})\|_{\mathcal{H}(k)} |\mathrm{d}P(\mathbf{x}) \tag{3}$$

$$= \int \sqrt{k(x,x)} dP(x) \|f\|_{\mathcal{H}(k)}$$
(4)

Lemma 6

If $\int \sqrt{k(\mathbf{x},\mathbf{x})} dP(\mathbf{x}) < \infty$ then $\mu_P(\mathbf{x}) \in \mathcal{H}(k)$.

Proof.

Consider the linear operator $Lf = \int f(x) dP(x)$ acting on $f \in \mathcal{H}(k)$. Claim that L is a bounded linear operator from $\mathcal{H}(k)$ to \mathbb{R} . Indeed,

$$|Lf| = \left| \int f(x) dP(x) \right| \le \int |f(x)| dP(x) \tag{1}$$

$$= \int |\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k)} | \mathrm{d}P(\mathbf{x})$$
 (2)

$$\leq \int \|f\|_{\mathcal{H}(k)} \|k(\cdot, x)\|_{\mathcal{H}(k)} |\mathrm{d}P(x) \tag{3}$$

$$= \int \sqrt{k(\mathbf{x}, \mathbf{x})} dP(\mathbf{x}) \|f\|_{\mathcal{H}(k)}$$
 (4)

Lemma 6

If $\int \sqrt{k(\mathbf{x}, \mathbf{x})} dP(\mathbf{x}) < \infty$ then $\mu_P(\mathbf{x}) \in \mathcal{H}(k)$.

Proof.

Consider the linear operator $Lf = \int f(x) dP(x)$ acting on $f \in \mathcal{H}(k)$. Claim that L is a bounded linear operator from $\mathcal{H}(k)$ to \mathbb{R} . Indeed,

$$|Lf| = \left| \int f(x) dP(x) \right| \le \int |f(x)| dP(x) \tag{1}$$

$$= \int |\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k)} | \mathrm{d}P(\mathbf{x})$$
 (2)

$$\leq \int \|f\|_{\mathcal{H}(k)} \|k(\cdot, \mathbf{x})\|_{\mathcal{H}(k)} |\mathrm{d}P(\mathbf{x}) \tag{3}$$

$$= \int \sqrt{k(\mathbf{x}, \mathbf{x})} dP(\mathbf{x}) \|f\|_{\mathcal{H}(k)}$$
 (4)

Lemma 6

If $\int \sqrt{k(\mathbf{x},\mathbf{x})} dP(\mathbf{x}) < \infty$ then $\mu_P(\mathbf{x}) \in \mathcal{H}(k)$.

Proof.

Consider the linear operator $Lf = \int f(x) dP(x)$ acting on $f \in \mathcal{H}(k)$. Claim that L is a bounded linear operator from $\mathcal{H}(k)$ to \mathbb{R} . Indeed,

$$|Lf| = \left| \int f(\mathbf{x}) dP(\mathbf{x}) \right| \le \int |f(\mathbf{x})| dP(\mathbf{x})$$
 (1)

$$= \int |\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k)} | \mathrm{d}P(\mathbf{x})$$
 (2)

$$\leq \int \|f\|_{\mathcal{H}(k)} \|k(\cdot, x)\|_{\mathcal{H}(k)} |\mathrm{d}P(x) \tag{3}$$

$$= \int \sqrt{k(\mathbf{x}, \mathbf{x})} dP(\mathbf{x}) \|f\|_{\mathcal{H}(k)}$$
 (4)

Lemma 6

If $\int \sqrt{k(\mathbf{x}, \mathbf{x})} dP(\mathbf{x}) < \infty$ then $\mu_P(\mathbf{x}) \in \mathcal{H}(k)$.

Proof.

Consider the linear operator $Lf = \int f(\mathbf{x}) dP(\mathbf{x})$ acting on $f \in \mathcal{H}(k)$. Claim that L is a bounded linear operator from $\mathcal{H}(k)$ to \mathbb{R} . Indeed,

$$|Lf| = \left| \int f(\mathbf{x}) dP(\mathbf{x}) \right| \le \int |f(\mathbf{x})| dP(\mathbf{x})$$
(1)

$$= \int |\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k)} | dP(\mathbf{x})$$
 (2)

$$\leq \int \|f\|_{\mathcal{H}(k)} \|k(\cdot, \mathbf{x})\|_{\mathcal{H}(k)} |\mathrm{d}P(\mathbf{x}) \tag{3}$$

$$= \int \sqrt{k(\mathbf{x}, \mathbf{x})} dP(\mathbf{x}) \|f\|_{\mathcal{H}(k)}$$
 (4)

Lemma 6

If $\int \sqrt{k(\mathbf{x},\mathbf{x})} \mathrm{d}P(\mathbf{x}) < \infty$ then $\mu_P(\mathbf{x}) \in \mathcal{H}(k)$.

Proof.

Consider the linear operator $Lf = \int f(x) dP(x)$ acting on $f \in \mathcal{H}(k)$. Claim that L is a bounded linear operator from $\mathcal{H}(k)$ to \mathbb{R} . Indeed,

$$|Lf| = \left| \int f(\mathbf{x}) dP(\mathbf{x}) \right| \le \int |f(\mathbf{x})| dP(\mathbf{x})$$
 (1)

$$= \int |\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k)} | dP(\mathbf{x})$$
 (2)

$$\leq \int \|f\|_{\mathcal{H}(k)} \|k(\cdot, \mathbf{x})\|_{\mathcal{H}(k)} |\mathrm{d}P(\mathbf{x}) \tag{3}$$

$$= \int \sqrt{k(\mathbf{x}, \mathbf{x})} dP(\mathbf{x}) \|f\|_{\mathcal{H}(k)}$$
(4)



Lemma 6

If $\int \sqrt{k(\mathbf{x},\mathbf{x})} \mathrm{d}P(\mathbf{x}) < \infty$ then $\mu_P(\mathbf{x}) \in \mathcal{H}(k)$.

Proof.

Consider the linear operator $Lf = \int f(x) dP(x)$ acting on $f \in \mathcal{H}(k)$. Claim that L is a bounded linear operator from $\mathcal{H}(k)$ to \mathbb{R} . Indeed,

$$|Lf| = \left| \int f(\mathbf{x}) dP(\mathbf{x}) \right| \le \int |f(\mathbf{x})| dP(\mathbf{x})$$
 (1)

$$= \int |\langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k)} | dP(\mathbf{x})$$
 (2)

$$\leq \int \|f\|_{\mathcal{H}(k)} \|k(\cdot, \mathbf{x})\|_{\mathcal{H}(k)} |\mathrm{d}P(\mathbf{x}) \tag{3}$$

$$= \int \sqrt{k(\mathbf{x}, \mathbf{x})} \mathrm{d}P(\mathbf{x}) \|f\|_{\mathcal{H}(k)} \tag{4}$$

Proof of Lemma 6, continued.

Thus, from the Riesz representation theorem, there exists $h \in \mathcal{H}(k)$ such that

$$Lf = \langle f, h \rangle_{\mathcal{H}(k)}.$$

Taking f(x) = k(y, x) and using the reproducing property leads to

$$\int k(y,x)dP(x) = Lf$$

$$= \langle f, h \rangle_{\mathcal{H}(k)}$$

$$= h(y),$$

so that $h(\cdot) = \int k(\cdot, \mathbf{x}) dP(\mathbf{x})$, and so $\mu_p = h \in \mathcal{H}(k)$, as was claimed.

Standing Assumption 1

Proof of Lemma 6, continued.

Thus, from the Riesz representation theorem, there exists $h \in \mathcal{H}(k)$ such that

$$Lf = \langle f, h \rangle_{\mathcal{H}(k)}.$$

Taking f(x) = k(y, x) and using the reproducing property leads to

$$\int k(y, x) dP(x) = Lf$$

$$= \langle f, h \rangle_{\mathcal{H}(k)}$$

$$= h(y),$$

so that $h(\cdot)=\int k(\cdot,x)\mathrm{d}P(x)$, and so $\mu_p=h\in\mathcal{H}(k)$, as was claimed.

Standing Assumption 1



Proof of Lemma 6, continued.

Thus, from the Riesz representation theorem, there exists $h \in \mathcal{H}(k)$ such that

$$Lf = \langle f, h \rangle_{\mathcal{H}(k)}.$$

Taking f(x) = k(y, x) and using the reproducing property leads to

$$\int k(y, x) dP(x) = Lf$$

$$= \langle f, h \rangle_{\mathcal{H}(k)}$$

$$= h(y),$$

so that $h(\cdot)=\int k(\cdot,x)\mathrm{d}P(x)$, and so $\mu_p=h\in\mathcal{H}(k)$, as was claimed.

Standing Assumption 1



Proof of Lemma 6, continued.

Thus, from the Riesz representation theorem, there exists $h \in \mathcal{H}(k)$ such that

$$Lf = \langle f, h \rangle_{\mathcal{H}(k)}.$$

Taking f(x) = k(y, x) and using the reproducing property leads to

$$\int k(\mathbf{y}, \mathbf{x}) dP(\mathbf{x}) = Lf$$

$$= \langle f, h \rangle_{\mathcal{H}(k)}$$

$$= h(\mathbf{y}),$$

so that $h(\cdot)=\int k(\cdot,x)\mathrm{d}P(x)$, and so $\mu_p=h\in\mathcal{H}(k)$, as was claimed.

Standing Assumption 1



Proof of Lemma 6, continued.

Thus, from the Riesz representation theorem, there exists $h \in \mathcal{H}(k)$ such that

$$Lf = \langle f, h \rangle_{\mathcal{H}(k)}.$$

Taking f(x) = k(y, x) and using the reproducing property leads to

$$\int k(\mathbf{y}, \mathbf{x}) dP(\mathbf{x}) = Lf$$

$$= \langle f, h \rangle_{\mathcal{H}(k)}$$

$$= h(\mathbf{y}),$$

so that $h(\cdot) = \int k(\cdot, \mathbf{x}) dP(\mathbf{x})$, and so $\mu_p = h \in \mathcal{H}(k)$, as was claimed.

Standing Assumption 1



Proof of Lemma 6, continued.

Thus, from the Riesz representation theorem, there exists $h \in \mathcal{H}(k)$ such that

$$Lf = \langle f, h \rangle_{\mathcal{H}(k)}.$$

Taking f(x) = k(y, x) and using the reproducing property leads to

$$\int k(y, x) dP(x) = Lf$$

$$= \langle f, h \rangle_{\mathcal{H}(k)}$$

$$= h(y),$$

so that $h(\cdot) = \int k(\cdot, \mathbf{x}) dP(\mathbf{x})$, and so $\mu_p = h \in \mathcal{H}(k)$, as was claimed.

Standing Assumption 1

For all reproducing kernels k and probability distributions P considered in the sequel, we assume that $\int \sqrt{k(\mathbf{x},\mathbf{x})} \mathrm{d}P(\mathbf{x}) < \infty$.

In this more general setting, the cubature error representer is the difference $e = \mu_{Q_n} - \mu_P$ of two kernel mean embeddings, where $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ is the discrete distribution on which the cubature rule is based. i.e.

$$\int f dP - \sum_{i=1}^n w_i f(\mathbf{x}_i) = \langle f, \mu_P - \mu_{Q_n} \rangle_{\mathcal{H}(k)}, \quad \mu_P = \int k(\cdot, \mathbf{x}) dP(\mathbf{x}), \quad \mu_{Q_n} = \sum_{i=1}^n w_i k(\cdot, \mathbf{x}_i).$$

There are several different ways to systematically assess the performance of a cubature rule, but here we focus on a worst case assessment:

Definition 7 (Maximum mean discrepancy)

The maximum mean discrepancy (MMD) between two distributions P and Q is

$$D_k(P,Q) = \|\mu_P - \mu_Q\|_{\mathcal{H}(k)} = \sup_{\|f\|_{\mathcal{H}(k)} \le 1} \left| \int f dP - \int f dQ \right|$$

In this more general setting, the cubature error representer is the difference $e=\mu_{Q_n}-\mu_P$ of two kernel mean embeddings, where $Q_n=\sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ is the discrete distribution on which the cubature rule is based. i.e.

$$\int f dP - \sum_{i=1}^n w_i f(\mathbf{x}_i) = \langle f, \mu_P - \mu_{Q_n} \rangle_{\mathcal{H}(k)}, \quad \mu_P = \int k(\cdot, \mathbf{x}) dP(\mathbf{x}), \quad \mu_{Q_n} = \sum_{i=1}^n w_i k(\cdot, \mathbf{x}_i).$$

There are several different ways to systematically assess the performance of a cubature rule, but here we focus on a *worst case* assessment:

Definition 7 (Maximum mean discrepancy)

The MMD between two distributions P and Q is

$$D_k(P,Q) = \|\mu_P - \mu_Q\|_{\mathcal{H}(k)} = \sup_{\|f\|_{\mathcal{H}(k)} \le 1} \left| \int f dP - \int f dQ \right|$$

In this more general setting, the cubature error representer is the difference $e = \mu_{Q_n} - \mu_P$ of two kernel mean embeddings, where $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ is the discrete distribution on which the cubature rule is based. i.e.

$$\int f dP - \sum_{i=1}^n w_i f(\mathbf{x}_i) = \langle f, \mu_P - \mu_{Q_n} \rangle_{\mathcal{H}(k)}, \quad \mu_P = \int k(\cdot, \mathbf{x}) dP(\mathbf{x}), \quad \mu_{Q_n} = \sum_{i=1}^n w_i k(\cdot, \mathbf{x}_i).$$

There are several different ways to systematically assess the performance of a cubature rule, but here we focus on a *worst case* assessment:

Definition 7 (Maximum mean discrepancy)

The MMD between two distributions P and Q is

$$D_k(P,Q) = \|\mu_P - \mu_Q\|_{\mathcal{H}(k)} = \sup_{\|f\|_{\mathcal{H}(k)} \leq 1} \left| \int f dP - \int f dQ \right|,$$

In this more general setting, the cubature error representer is the difference $e = \mu_{Q_n} - \mu_P$ of two kernel mean embeddings, where $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ is the discrete distribution on which the cubature rule is based. i.e.

$$\int f dP - \sum_{i=1}^n w_i f(\mathbf{x}_i) = \langle f, \mu_P - \mu_{Q_n} \rangle_{\mathcal{H}(k)}, \quad \mu_P = \int k(\cdot, \mathbf{x}) dP(\mathbf{x}), \quad \mu_{Q_n} = \sum_{i=1}^n w_i k(\cdot, \mathbf{x}_i).$$

There are several different ways to systematically assess the performance of a cubature rule, but here we focus on a *worst case* assessment:

Definition 7 (Maximum mean discrepancy)

The MMD between two distributions P and Q is

$$D_k(P,Q) = \|\mu_P - \mu_Q\|_{\mathcal{H}(k)} = \sup_{\|f\|_{\mathcal{H}(k)} \leq 1} \left| \int f dP - \int f dQ \right|,$$

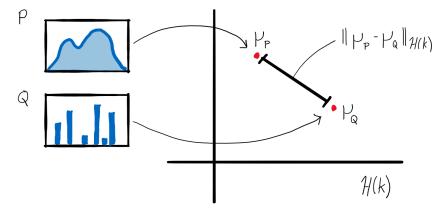


Figure: Kernel mean embedding: Two probability distributions P and Q are mapped to their respective elements μ_P and μ_Q in the RKHS $\mathcal{H}(k)$. The distance (in $\mathcal{H}(k)$) between these kernel mean embeddings μ_P and μ_Q is called the MMD between P and Q.

If $D_k(P, Q_n) = 0$, the cubature rule based on Q_n will be exact for all integrands $f \in \mathcal{H}(k)$. Does this mean that Q_n and P are identical?

Definition 8 (Characteristic kernel)

A kernel k is said to be *characteristic* if $D_k(P,Q) = 0$ implies P = Q.

Example 9 (Polynomial kernel is not characteristic)

From Example 2, the kernel $k(x,y) = \sum_{i=0}^{p} x^{i} y^{i}$ reproduces an RKHS whose elements are the polynomials of degree at most p on the domain $\mathcal{X} = \mathbb{R}$. Thus $D_k(P,Q) = 0$ if and only if the moments $\int x^{i} dP(x)$ and $\int x^{i} dQ(x)$ are identical for $i = 1, \ldots, p$. In particular, k is *not* a characteristic kernel.

Example 10

The Gaussian kernel $k(x,y) = \exp(-\|x-y\|^2)$ is a characteristic kernel on $\mathcal{X} = \mathbb{R}^d$



If $D_k(P, Q_n) = 0$, the cubature rule based on Q_n will be exact for all integrands $f \in \mathcal{H}(k)$. Does this mean that Q_n and P are identical?

Definition 8 (Characteristic kernel)

A kernel k is said to be *characteristic* if $D_k(P,Q) = 0$ implies P = Q.

Example 9 (Polynomial kernel is not characteristic)

From Example 2, the kernel $k(x,y) = \sum_{i=0}^{p} x^i y^i$ reproduces an RKHS whose elements are the polynomials of degree at most p on the domain $\mathcal{X} = \mathbb{R}$. Thus $D_k(P,Q) = 0$ if and only if the moments $\int x^i \mathrm{d}P(x)$ and $\int x^i \mathrm{d}Q(x)$ are identical for $i=1,\ldots,p$. In particular, k is not a characteristic kernel.

Example 10

The Gaussian kernel $k(x,y) = \exp(-\|x-y\|^2)$ is a characteristic kernel on $\mathcal{X} = \mathbb{R}^d$



If $D_k(P, Q_n) = 0$, the cubature rule based on Q_n will be exact for all integrands $f \in \mathcal{H}(k)$. Does this mean that Q_n and P are identical?

Definition 8 (Characteristic kernel)

A kernel k is said to be *characteristic* if $D_k(P,Q) = 0$ implies P = Q.

Example 9 (Polynomial kernel is not characteristic)

From Example 2, the kernel $k(x,y) = \sum_{i=0}^{p} x^i y^i$ reproduces an RKHS whose elements are the polynomials of degree at most p on the domain $\mathcal{X} = \mathbb{R}$. Thus $D_k(P,Q) = 0$ if and only if the moments $\int x^i \mathrm{d}P(x)$ and $\int x^i \mathrm{d}Q(x)$ are identical for $i=1,\ldots,p$. In particular, k is not a characteristic kernel.

Example 10

The Gaussian kernel $k(x,y) = \exp(-\|x-y\|^2)$ is a characteristic kernel on $\mathcal{X} = \mathbb{R}^d$.



If $D_k(P, Q_n) = 0$, the cubature rule based on Q_n will be exact for all integrands $f \in \mathcal{H}(k)$. Does this mean that Q_n and P are identical?

Definition 8 (Characteristic kernel)

A kernel k is said to be *characteristic* if $D_k(P,Q) = 0$ implies P = Q.

Example 9 (Polynomial kernel is not characteristic)

From Example 2, the kernel $k(x,y) = \sum_{i=0}^{p} x^i y^i$ reproduces an RKHS whose elements are the polynomials of degree at most p on the domain $\mathcal{X} = \mathbb{R}$. Thus $D_k(P,Q) = 0$ if and only if the moments $\int x^i \mathrm{d}P(x)$ and $\int x^i \mathrm{d}Q(x)$ are identical for $i=1,\ldots,p$. In particular, k is not a characteristic kernel.

Example 10

The Gaussian kernel $k(x, y) = \exp(-\|x - y\|^2)$ is a characteristic kernel on $\mathcal{X} = \mathbb{R}^d$.



Notation: Let $Q_n \Rightarrow P$ denote that the sequence $(Q_n)_{n=1}^{\infty}$ converges *weakly* (or *in distribution*) to P (i.e. $\int f dQ_n \to \int f dP$ for all functions f which are continuous and bounded).

Definition 11 (Weak convergence control)

A kernel k is said to have weak convergence control if $D_k(P, Q_n) \to 0$ implies that $Q_n \Rightarrow P$.

Remark 1

For a compact Hausdorff space \mathcal{X} , a bounded characteristic kernel k has weak convergence control. This need not hold when the domain \mathcal{X} is non-compact.

Convergence control justifies attempting to minimise MMD for the purposes of quantisation and more general approximation, as we will attempt in the sequel.

Example 12

Notation: Let $Q_n \Rightarrow P$ denote that the sequence $(Q_n)_{n=1}^{\infty}$ converges *weakly* (or *in distribution*) to P (i.e. $\int f dQ_n \to \int f dP$ for all functions f which are continuous and bounded).

Definition 11 (Weak convergence control)

A kernel k is said to have weak convergence control if $D_k(P, Q_n) \to 0$ implies that $Q_n \Rightarrow P$.

Remark 1

For a compact Hausdorff space \mathcal{X} , a bounded characteristic kernel k has weak convergence control. This need not hold when the domain \mathcal{X} is non-compact.

Convergence control justifies attempting to minimise MMD for the purposes of quantisation and more general approximation, as we will attempt in the sequel.

Example 12

Notation: Let $Q_n \Rightarrow P$ denote that the sequence $(Q_n)_{n=1}^{\infty}$ converges *weakly* (or *in distribution*) to P (i.e. $\int f dQ_n \to \int f dP$ for all functions f which are continuous and bounded).

Definition 11 (Weak convergence control)

A kernel k is said to have weak convergence control if $D_k(P, Q_n) \to 0$ implies that $Q_n \Rightarrow P$.

Remark 1

For a compact Hausdorff space \mathcal{X} , a bounded characteristic kernel k has weak convergence control. This need not hold when the domain \mathcal{X} is non-compact.

Convergence control justifies attempting to minimise MMD for the purposes of quantisation and more general approximation, as we will attempt in the sequel.

Example 12

Notation: Let $Q_n \Rightarrow P$ denote that the sequence $(Q_n)_{n=1}^{\infty}$ converges *weakly* (or *in distribution*) to P (i.e. $\int f dQ_n \to \int f dP$ for all functions f which are continuous and bounded).

Definition 11 (Weak convergence control)

A kernel k is said to have weak convergence control if $D_k(P, Q_n) \to 0$ implies that $Q_n \Rightarrow P$.

Remark 1

For a compact Hausdorff space \mathcal{X} , a bounded characteristic kernel k has weak convergence control. This need not hold when the domain \mathcal{X} is non-compact.

Convergence control justifies attempting to minimise MMD for the purposes of quantisation and more general approximation, as we will attempt in the sequel.

Example 12

Notation: Let $Q_n \Rightarrow P$ denote that the sequence $(Q_n)_{n=1}^{\infty}$ converges *weakly* (or *in distribution*) to P (i.e. $\int f dQ_n \to \int f dP$ for all functions f which are continuous and bounded).

Definition 11 (Weak convergence control)

A kernel k is said to have weak convergence control if $D_k(P, Q_n) \to 0$ implies that $Q_n \Rightarrow P$.

Remark 1

For a compact Hausdorff space \mathcal{X} , a bounded characteristic kernel k has weak convergence control. This need not hold when the domain \mathcal{X} is non-compact.

Convergence control justifies attempting to minimise MMD for the purposes of quantisation and more general approximation, as we will attempt in the sequel.

Example 12

MMD is a convenient measure of sample quality because it can be computed:

$$\begin{aligned} D_k(P,Q)^2 &= \|\mu_P - \mu_Q\|_{\mathcal{H}(k)}^2 \\ &= \langle \mu_P - \mu_Q, \mu_P - \mu_Q \rangle_{\mathcal{H}(k)} \\ &= \langle \mu_P, \mu_P \rangle_{\mathcal{H}(k)} - 2\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)} + \langle \mu_Q, \mu_Q \rangle_{\mathcal{H}(k)} \end{aligned}$$

Considering for example the term $\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)}$, we have

$$\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} = \left\langle \int k(\cdot, x) dP(x), \int k(\cdot, y) dQ(y) \right\rangle_{\mathcal{H}(k)}$$
$$= \iint \langle k(\cdot, x), k(\cdot, y) \rangle_{\mathcal{H}(k)} dP(x) dQ(y) = \iint k(x, y) dP(x) dQ(y).$$

Here we have used the reproducing property, as well as using Lemma 6 to justify the exchanges of integral and inner product. Proceeding similarly with all three terms results in the expression

$$D_k(P,Q)^2 = \iint k(x,y) \mathrm{d}P(x) \mathrm{d}P(y) - 2 \iint k(x,y) \mathrm{d}P(x) \mathrm{d}Q(y) + \iint k(x,y) \mathrm{d}Q(x) \mathrm{d}Q(y).$$

How to exploit MMD for sampling?

MMD is a convenient measure of sample quality because it can be computed:

$$\begin{aligned} D_{k}(P,Q)^{2} &= \|\mu_{P} - \mu_{Q}\|_{\mathcal{H}(k)}^{2} \\ &= \langle \mu_{P} - \mu_{Q}, \mu_{P} - \mu_{Q} \rangle_{\mathcal{H}(k)} \\ &= \langle \mu_{P}, \mu_{P} \rangle_{\mathcal{H}(k)} - 2\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} + \langle \mu_{Q}, \mu_{Q} \rangle_{\mathcal{H}(k)} \end{aligned}$$

Considering for example the term $\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)}$, we have

$$\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} = \left\langle \int k(\cdot, x) dP(x), \int k(\cdot, y) dQ(y) \right\rangle_{\mathcal{H}(k)}$$
$$= \iint \langle k(\cdot, x), k(\cdot, y) \rangle_{\mathcal{H}(k)} dP(x) dQ(y) = \iint k(x, y) dP(x) dQ(y).$$

Here we have used the reproducing property, as well as using Lemma 6 to justify the exchanges of integral and inner product. Proceeding similarly with all three terms results in the expression

$$D_k(P,Q)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y).$$

MMD is a convenient measure of sample quality because it can be computed:

$$\begin{split} D_{k}(P,Q)^{2} &= \|\mu_{P} - \mu_{Q}\|_{\mathcal{H}(k)}^{2} \\ &= \langle \mu_{P} - \mu_{Q}, \mu_{P} - \mu_{Q} \rangle_{\mathcal{H}(k)} \\ &= \langle \mu_{P}, \mu_{P} \rangle_{\mathcal{H}(k)} - 2\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} + \langle \mu_{Q}, \mu_{Q} \rangle_{\mathcal{H}(k)} \end{split}$$

Considering for example the term $\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)}$, we have

$$\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} = \left\langle \int k(\cdot, x) dP(x), \int k(\cdot, y) dQ(y) \right\rangle_{\mathcal{H}(k)}$$
$$= \iint \langle k(\cdot, x), k(\cdot, y) \rangle_{\mathcal{H}(k)} dP(x) dQ(y) = \iint k(x, y) dP(x) dQ(y).$$

Here we have used the reproducing property, as well as using Lemma 6 to justify the exchanges of integral and inner product. Proceeding similarly with all three terms results in the expression

$$D_k(P,Q)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y).$$

MMD is a convenient measure of sample quality because it can be computed:

$$\begin{split} D_k(P,Q)^2 &= \|\mu_P - \mu_Q\|_{\mathcal{H}(k)}^2 \\ &= \langle \mu_P - \mu_Q, \mu_P - \mu_Q \rangle_{\mathcal{H}(k)} \\ &= \langle \mu_P, \mu_P \rangle_{\mathcal{H}(k)} - 2\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)} + \langle \mu_Q, \mu_Q \rangle_{\mathcal{H}(k)}. \end{split}$$

Considering for example the term $\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)}$, we have

$$\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} = \left\langle \int k(\cdot, x) dP(x), \int k(\cdot, y) dQ(y) \right\rangle_{\mathcal{H}(k)}$$
$$= \iint \langle k(\cdot, x), k(\cdot, y) \rangle_{\mathcal{H}(k)} dP(x) dQ(y) = \iint k(x, y) dP(x) dQ(y).$$

Here we have used the reproducing property, as well as using Lemma 6 to justify the exchanges of integral and inner product. Proceeding similarly with all three terms results in the expression

$$D_k(P,Q)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y).$$

MMD is a convenient measure of sample quality because it can be computed:

$$\begin{split} D_k(P,Q)^2 &= \|\mu_P - \mu_Q\|_{\mathcal{H}(k)}^2 \\ &= \langle \mu_P - \mu_Q, \mu_P - \mu_Q \rangle_{\mathcal{H}(k)} \\ &= \langle \mu_P, \mu_P \rangle_{\mathcal{H}(k)} - 2\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)} + \langle \mu_Q, \mu_Q \rangle_{\mathcal{H}(k)}. \end{split}$$

Considering for example the term $\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)}$, we have

$$\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} = \left\langle \int k(\cdot, x) dP(x), \int k(\cdot, y) dQ(y) \right\rangle_{\mathcal{H}(k)}$$
$$= \iint \langle k(\cdot, x), k(\cdot, y) \rangle_{\mathcal{H}(k)} dP(x) dQ(y) = \iint k(x, y) dP(x) dQ(y).$$

Here we have used the reproducing property, as well as using Lemma 6 to justify the exchanges of integral and inner product. Proceeding similarly with all three terms results in the expression

$$D_k(P,Q)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y).$$

MMD is a convenient measure of sample quality because it can be computed:

$$\begin{split} D_k(P,Q)^2 &= \|\mu_P - \mu_Q\|_{\mathcal{H}(k)}^2 \\ &= \langle \mu_P - \mu_Q, \mu_P - \mu_Q \rangle_{\mathcal{H}(k)} \\ &= \langle \mu_P, \mu_P \rangle_{\mathcal{H}(k)} - 2\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)} + \langle \mu_Q, \mu_Q \rangle_{\mathcal{H}(k)}. \end{split}$$

Considering for example the term $\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)}$, we have

$$\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} = \left\langle \int k(\cdot, \mathbf{x}) dP(\mathbf{x}), \int k(\cdot, \mathbf{y}) dQ(\mathbf{y}) \right\rangle_{\mathcal{H}(k)}$$
$$= \int \int \langle k(\cdot, \mathbf{x}), k(\cdot, \mathbf{y}) \rangle_{\mathcal{H}(k)} dP(\mathbf{x}) dQ(\mathbf{y}) = \int \int k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) dQ(\mathbf{y}).$$

Here we have used the reproducing property, as well as using Lemma 6 to justify the exchanges of integral and inner product. Proceeding similarly with all three terms results in the expression

$$D_k(P,Q)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y).$$

MMD is a convenient measure of sample quality because it can be computed:

$$\begin{split} D_k(P,Q)^2 &= \|\mu_P - \mu_Q\|_{\mathcal{H}(k)}^2 \\ &= \langle \mu_P - \mu_Q, \mu_P - \mu_Q \rangle_{\mathcal{H}(k)} \\ &= \langle \mu_P, \mu_P \rangle_{\mathcal{H}(k)} - 2\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)} + \langle \mu_Q, \mu_Q \rangle_{\mathcal{H}(k)}. \end{split}$$

Considering for example the term $\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)}$, we have

$$\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} = \left\langle \int k(\cdot, \mathbf{x}) dP(\mathbf{x}), \int k(\cdot, \mathbf{y}) dQ(\mathbf{y}) \right\rangle_{\mathcal{H}(k)}$$
$$= \iint \langle k(\cdot, \mathbf{x}), k(\cdot, \mathbf{y}) \rangle_{\mathcal{H}(k)} dP(\mathbf{x}) dQ(\mathbf{y}) = \iint k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) dQ(\mathbf{y}).$$

Here we have used the reproducing property, as well as using Lemma 6 to justify the exchanges of integral and inner product. Proceeding similarly with all three terms results in the expression

$$D_k(P,Q)^2 = \iint k(x,y) \mathrm{d}P(x) \mathrm{d}P(y) - 2 \iint k(x,y) \mathrm{d}P(x) \mathrm{d}Q(y) + \iint k(x,y) \mathrm{d}Q(x) \mathrm{d}Q(y).$$



MMD is a convenient measure of sample quality because it can be computed:

$$\begin{split} D_k(P,Q)^2 &= \|\mu_P - \mu_Q\|_{\mathcal{H}(k)}^2 \\ &= \langle \mu_P - \mu_Q, \mu_P - \mu_Q \rangle_{\mathcal{H}(k)} \\ &= \langle \mu_P, \mu_P \rangle_{\mathcal{H}(k)} - 2\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)} + \langle \mu_Q, \mu_Q \rangle_{\mathcal{H}(k)}. \end{split}$$

Considering for example the term $\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)}$, we have

$$\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} = \left\langle \int k(\cdot, \mathbf{x}) dP(\mathbf{x}), \int k(\cdot, \mathbf{y}) dQ(\mathbf{y}) \right\rangle_{\mathcal{H}(k)}$$
$$= \iint \langle k(\cdot, \mathbf{x}), k(\cdot, \mathbf{y}) \rangle_{\mathcal{H}(k)} dP(\mathbf{x}) dQ(\mathbf{y}) = \iint k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) dQ(\mathbf{y}).$$

Here we have used the reproducing property, as well as using Lemma 6 to justify the exchanges of integral and inner product. Proceeding similarly with all three terms results in the expression

$$D_k(P,Q)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y).$$



MMD is a convenient measure of sample quality because it can be computed:

$$\begin{split} D_k(P,Q)^2 &= \|\mu_P - \mu_Q\|_{\mathcal{H}(k)}^2 \\ &= \langle \mu_P - \mu_Q, \mu_P - \mu_Q \rangle_{\mathcal{H}(k)} \\ &= \langle \mu_P, \mu_P \rangle_{\mathcal{H}(k)} - 2\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)} + \langle \mu_Q, \mu_Q \rangle_{\mathcal{H}(k)}. \end{split}$$

Considering for example the term $\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)}$, we have

$$\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} = \left\langle \int k(\cdot, \mathbf{x}) dP(\mathbf{x}), \int k(\cdot, \mathbf{y}) dQ(\mathbf{y}) \right\rangle_{\mathcal{H}(k)}$$
$$= \iint \langle k(\cdot, \mathbf{x}), k(\cdot, \mathbf{y}) \rangle_{\mathcal{H}(k)} dP(\mathbf{x}) dQ(\mathbf{y}) = \iint k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) dQ(\mathbf{y}).$$

Here we have used the reproducing property, as well as using Lemma 6 to justify the exchanges of integral and inner product. Proceeding similarly with all three terms results in the expression

$$D_k(P,Q)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y).$$



MMD is a convenient measure of sample quality because it can be computed:

$$\begin{split} D_k(P,Q)^2 &= \|\mu_P - \mu_Q\|_{\mathcal{H}(k)}^2 \\ &= \langle \mu_P - \mu_Q, \mu_P - \mu_Q \rangle_{\mathcal{H}(k)} \\ &= \langle \mu_P, \mu_P \rangle_{\mathcal{H}(k)} - 2\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)} + \langle \mu_Q, \mu_Q \rangle_{\mathcal{H}(k)}. \end{split}$$

Considering for example the term $\langle \mu_P, \mu_Q \rangle_{\mathcal{H}(k)}$, we have

$$\langle \mu_{P}, \mu_{Q} \rangle_{\mathcal{H}(k)} = \left\langle \int k(\cdot, \mathbf{x}) dP(\mathbf{x}), \int k(\cdot, \mathbf{y}) dQ(\mathbf{y}) \right\rangle_{\mathcal{H}(k)}$$
$$= \iint \langle k(\cdot, \mathbf{x}), k(\cdot, \mathbf{y}) \rangle_{\mathcal{H}(k)} dP(\mathbf{x}) dQ(\mathbf{y}) = \iint k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) dQ(\mathbf{y}).$$

Here we have used the reproducing property, as well as using Lemma 6 to justify the exchanges of integral and inner product. Proceeding similarly with all three terms results in the expression

$$D_k(P,Q)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y).$$



Sampling with Kernels

The goal of *quantisation* is to find Q of the form $Q_n = \frac{1}{n} \sum_{i=1}^n \delta(\mathbf{x}_i)$ such that $Q_n \approx P$.

As a simple baseline method for quantisation we consider Monte Carlo:

Figure: Monte Carlo: Independent samples (blue circles) from a "horseshoe" distribution P (grey contours).

The goal of *quantisation* is to find Q of the form $Q_n = \frac{1}{n} \sum_{i=1}^n \delta(\mathbf{x}_i)$ such that $Q_n \approx P$.

As a simple baseline method for quantisation we consider Monte Carlo:

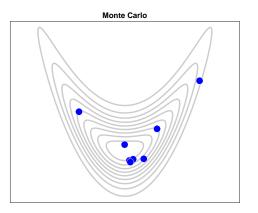


Figure: Monte Carlo: Independent samples (blue circles) from a "horseshoe" distribution P (grey contours).

Proposition 1 (MMD of Monte Carlo)

Let $x_1, \ldots, x_n \sim P$ be independent. Assume that $C := \int k(x, x) dP(x) < \infty$. Then

$$\mathbb{E}\left[D_k(P,Q_n)^2\right]\leq \frac{C}{n}.$$

Proof

From the closed form of MMD, with $Q = Q_n = \frac{1}{n} \sum_{i=1}^n \delta(x_i)$, we obtain that

$$D_k(P,Q_n)^2 = \iint k(x,y) dP(x) dP(y) - \frac{2}{n} \sum_{i=1}^n \int k(x,x_i) dP(x) + \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j)$$

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] = \mathbb{E}\left[\frac{1}{n^2}\sum_{i=1}^n\sum_{j=1}^n k(x_i,x_j)\right] - \iint k(x,y)\mathrm{d}P(x)\mathrm{d}P(y)$$

Proposition 1 (MMD of Monte Carlo)

Let $x_1, \ldots, x_n \sim P$ be independent. Assume that $C := \int k(x, x) dP(x) < \infty$. Then

$$\mathbb{E}\left[D_k(P,Q_n)^2\right]\leq \frac{C}{n}.$$

Proof.

From the closed form of MMD, with $Q = Q_n = \frac{1}{n} \sum_{i=1}^n \delta(\mathbf{x}_i)$, we obtain that

$$D_k(P,Q_n)^2 = \iint k(x,y) dP(x) dP(y) - \frac{2}{n} \sum_{i=1}^n \int k(x,x_i) dP(x) + \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j) dP(x) dP$$

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] = \mathbb{E}\left[\frac{1}{n^2}\sum_{i=1}^n\sum_{j=1}^n k(x_i,x_j)\right] - \iint k(x,y)\mathrm{d}P(x)\mathrm{d}P(y)$$



Proposition 1 (MMD of Monte Carlo)

Let $x_1, \ldots, x_n \sim P$ be independent. Assume that $C := \int k(x, x) dP(x) < \infty$. Then

$$\mathbb{E}\left[D_k(P,Q_n)^2\right]\leq \frac{C}{n}.$$

Proof.

From the closed form of MMD, with $Q = Q_n = \frac{1}{n} \sum_{i=1}^n \delta(\mathbf{x}_i)$, we obtain that

$$D_k(P,Q_n)^2 = \iint k(x,y) dP(x) dP(y) - \frac{2}{n} \sum_{i=1}^n \int k(x,x_i) dP(x) + \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j) dP(x) dP$$

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] = \mathbb{E}\left[\frac{1}{n^2}\sum_{i=1}^n\sum_{j=1}^n k(x_i,x_j)\right] - \iint k(x,y)dP(x)dP(y)$$



Proposition 1 (MMD of Monte Carlo)

Let $x_1, \ldots, x_n \sim P$ be independent. Assume that $C := \int k(x, x) dP(x) < \infty$. Then

$$\mathbb{E}\left[D_k(P,Q_n)^2\right]\leq \frac{C}{n}.$$

Proof.

From the closed form of MMD, with $Q = Q_n = \frac{1}{n} \sum_{i=1}^n \delta(\mathbf{x}_i)$, we obtain that

$$D_k(P,Q_n)^2 = \iint k(x,y) dP(x) dP(y) - \frac{2}{n} \sum_{i=1}^n \int k(x,x_i) dP(x) + \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j) dP(x) dP$$

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] = \mathbb{E}\left[\frac{1}{n^2}\sum_{i=1}^n k(\mathbf{x}_i,\mathbf{x}_i) + \frac{1}{n^2}\sum_{i=1}^n\sum_{j\neq i} k(\mathbf{x}_i,\mathbf{x}_j)\right] - \iint k(\mathbf{x},\mathbf{y})\mathrm{d}P(\mathbf{x})\mathrm{d}P(\mathbf{y})$$

Proposition 1 (MMD of Monte Carlo)

Let $x_1, \ldots, x_n \sim P$ be independent. Assume that $C := \int k(x, x) dP(x) < \infty$. Then

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] \leq \frac{C}{n}.$$

Proof.

From the closed form of MMD, with $Q = Q_n = \frac{1}{n} \sum_{i=1}^n \delta(\mathbf{x}_i)$, we obtain that

$$D_k(P,Q_n)^2 = \iint k(x,y) dP(x) dP(y) - \frac{2}{n} \sum_{i=1}^n \int k(x,x_i) dP(x) + \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j)$$

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] = \mathbb{E}\left[\frac{1}{n^2}\sum_{i=1}^n k(\mathbf{x}_i,\mathbf{x}_i)\right] - \underbrace{\frac{1}{n}\iint k(\mathbf{x},\mathbf{y})\mathrm{d}P(\mathbf{x})\mathrm{d}P(\mathbf{y})}_{>0}$$

Proposition 1 (MMD of Monte Carlo)

Let $x_1, \ldots, x_n \sim P$ be independent. Assume that $C := \int k(x, x) \mathrm{d}P(x) < \infty$. Then

$$\mathbb{E}\left[D_k(P,Q_n)^2\right]\leq \frac{C}{n}.$$

Proof.

From the closed form of MMD, with $Q = Q_n = \frac{1}{n} \sum_{i=1}^n \delta(\mathbf{x}_i)$, we obtain that

$$D_k(P,Q_n)^2 = \iint k(x,y) dP(x) dP(y) - \frac{2}{n} \sum_{i=1}^n \int k(x,x_i) dP(x) + \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j) dP(y) dP(y) = \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j) dP(y) + \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j) dP(y) dP(y) + \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j) dP(y) dP(y) + \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j) dP(y) dP(y) dP(y) + \frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n k(x_i,x_j) dP(y) dP($$

Taking expectations of both sides gives

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] \leq \frac{1}{n} \int k(\mathbf{x},\mathbf{x}) \mathrm{d}P(\mathbf{x})$$

as claimed.

Monte Carlo sampling provides a consistent but potentially far from optimal quantisation of P.

Note that the convergence rate in Proposition 1 does not depend on the kernel k, which highlights the inefficiency of the Monte Carlo method. [But the rate is dimension-independent.]

The goal of *optimal* quantisation is to approximate P as well as possible, with a fixed number of states x_1, \ldots, x_n .

In what follows we are going to introduce some algorithms that are **often inferior to QMC**, but have the advantage of being **applicable to general Bayesian inference tasks** with minor modification, as we will see in the sequel.

Algorithm 1 (Optimise Everything)

Approximate a solution to

$$\min_{(\mathbf{x}_1,\ldots,\mathbf{x}_n)\in\mathcal{X}\times\cdots\times\mathcal{X}} D_k\left(P,\frac{1}{n}\sum_{i=1}^n \delta(\mathbf{x}_i)\right)$$

Monte Carlo sampling provides a consistent but potentially far from optimal quantisation of P.

Note that the convergence rate in Proposition 1 does not depend on the kernel k, which highlights the inefficiency of the Monte Carlo method. [But the rate is dimension-independent.]

The goal of *optimal* quantisation is to approximate P as well as possible, with a fixed number of states x_1, \ldots, x_n .

In what follows we are going to introduce some algorithms that are **often inferior to QMC**, but have the advantage of being **applicable to general Bayesian inference tasks** with minor modification, as we will see in the sequel.

Algorithm 1 (Optimise Everything)

Approximate a solution to

$$\min_{(\mathbf{x}_1,\ldots,\mathbf{x}_n)\in\mathcal{X}\times\cdots\times\mathcal{X}} D_k\left(P,\frac{1}{n}\sum_{i=1}^n \delta(\mathbf{x}_i)\right)$$

Monte Carlo sampling provides a consistent but potentially far from optimal quantisation of P.

Note that the convergence rate in Proposition 1 does not depend on the kernel k, which highlights the inefficiency of the Monte Carlo method. [But the rate is dimension-independent.]

The goal of *optimal* quantisation is to approximate P as well as possible, with a fixed number of states x_1, \ldots, x_n .

In what follows we are going to introduce some algorithms that are **often inferior to QMC**, but have the advantage of being **applicable to general Bayesian inference tasks** with minor modification, as we will see in the sequel.

Algorithm 1 (Optimise Everything)

Approximate a solution to

$$\min_{(x_1,\ldots,x_n)\in\mathcal{X}\times\cdots\times\mathcal{X}} D_k\left(P,\frac{1}{n}\sum_{i=1}^n \delta(x_i)\right)$$

Monte Carlo sampling provides a consistent but potentially far from optimal quantisation of P.

Note that the convergence rate in Proposition 1 does not depend on the kernel k, which highlights the inefficiency of the Monte Carlo method. [But the rate is dimension-independent.]

The goal of *optimal* quantisation is to approximate P as well as possible, with a fixed number of states x_1, \ldots, x_n .

In what follows we are going to introduce some algorithms that are **often inferior to QMC**, but have the advantage of being **applicable to general Bayesian inference tasks** with minor modification, as we will see in the sequel.

Algorithm 1 (Optimise Everything)

Approximate a solution to

$$\min_{(x_1,\ldots,x_n)\in\mathcal{X}\times\cdots\times\mathcal{X}} D_k\left(P,\frac{1}{n}\sum_{i=1}^n \delta(x_i)\right)$$

Monte Carlo sampling provides a consistent but potentially far from optimal quantisation of P.

Note that the convergence rate in Proposition 1 does not depend on the kernel k, which highlights the inefficiency of the Monte Carlo method. [But the rate is dimension-independent.]

The goal of *optimal* quantisation is to approximate P as well as possible, with a fixed number of states x_1, \ldots, x_n .

In what follows we are going to introduce some algorithms that are **often inferior to QMC**, but have the advantage of being **applicable to general Bayesian inference tasks** with minor modification, as we will see in the sequel.

Algorithm 1 (Optimise Everything

Approximate a solution to

$$\min_{(x_1,\ldots,x_n)\in\mathcal{X}\times\cdots\times\mathcal{X}} D_k\left(P,\frac{1}{n}\sum_{i=1}^n \delta(x_i)\right)$$

Monte Carlo sampling provides a consistent but potentially far from optimal quantisation of P.

Note that the convergence rate in Proposition 1 does not depend on the kernel k, which highlights the inefficiency of the Monte Carlo method. [But the rate is dimension-independent.]

The goal of *optimal* quantisation is to approximate P as well as possible, with a fixed number of states x_1, \ldots, x_n .

In what follows we are going to introduce some algorithms that are **often inferior to QMC**, but have the advantage of being **applicable to general Bayesian inference tasks** with minor modification, as we will see in the sequel.

Algorithm 1 (Optimise Everything)

Approximate a solution to

$$\min_{(\mathbf{x}_1,\dots,\mathbf{x}_n)\in\mathcal{X}\times\dots\times\mathcal{X}} D_k\left(P,\frac{1}{n}\sum_{i=1}^n \delta(\mathbf{x}_i)\right)$$







Figure: Optimal quantisation: Global (continuous) minimisation of MMD to select states $\{x_i\}_{i=1}^n$ in an approximation Q_n to P. Example due to Gräf et al. [2012].

This example is the *stippling* procedure discussed in Gräf et al. [2012], implemented using a bandlimited kernel and the nonlinear CG method. See also *minimum energy designs* and MMD *gradient flow* [Arbel et al., 2019]. Requires a good initialisation to work well.

Also requires selecting n in advance - a non-extensible sequence of approximations to P





Figure: Optimal quantisation: Global (continuous) minimisation of MMD to select states $\{x_i\}_{i=1}^n$ in an approximation Q_n to P. Example due to Gräf et al. [2012].

This example is the *stippling* procedure discussed in Gräf et al. [2012], implemented using a bandlimited kernel and the nonlinear CG method. See also *minimum energy designs* and MMD *gradient flow* [Arbel et al., 2019]. Requires a good initialisation to work well.

Also requires selecting n in advance - a non-extensible sequence of approximations to P





Figure: Optimal quantisation: Global (continuous) minimisation of MMD to select states $\{x_i\}_{i=1}^n$ in an approximation Q_n to P. Example due to Gräf et al. [2012].

This example is the *stippling* procedure discussed in Gräf et al. [2012], implemented using a bandlimited kernel and the nonlinear CG method. See also *minimum energy designs* and MMD gradient flow [Arbel et al., 2019]. Requires a good initialisation to work well.

Also requires selecting n in advance - a non-extensible sequence of approximations to P.

To arrive at an extensible approximation, consider a greedy sequential algorithm:

Algorithm 2 (Continuous Greedy Optimisation)

Approximate a solution to

$$x_n \in \underset{x \in \mathcal{X}}{\operatorname{arg\,min}} \ D_k\left(P, \frac{1}{n}\delta(x) + \frac{1}{n}\sum_{i=1}^{n-1}\delta(x_i)\right), \qquad n = 1, 2, \dots$$

To arrive at an extensible approximation, consider a greedy sequential algorithm:

Algorithm 2 (Continuous Greedy Optimisation)

Approximate a solution to

$$\mathbf{x}_n \in \operatorname*{arg\,min}_{\mathbf{x} \in \mathcal{X}} \ D_k \left(P, \frac{1}{n} \delta(\mathbf{x}) + \frac{1}{n} \sum_{i=1}^{n-1} \delta(\mathbf{x}_i) \right), \qquad n = 1, 2, \dots$$

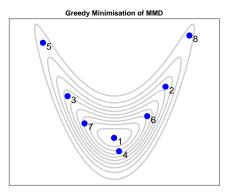


Figure: Optimal quantisation: Sequential (greedy) minimisation of MMD to select n = 8 states $\{x_i\}_{i=1}^n$ in an approximation Q_n to P. The numbers indicate the order in which the states x_i were selected.

However, the computation required to select the location x_n becomes increasingly difficult as n is increased. A computationally simpler algorithm is proposed next.

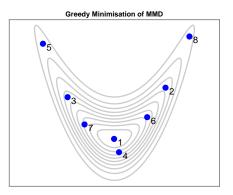


Figure: Optimal quantisation: Sequential (greedy) minimisation of MMD to select n = 8 states $\{x_i\}_{i=1}^n$ in an approximation Q_n to P. The numbers indicate the order in which the states x_i were selected.

However, the computation required to select the location x_n becomes increasingly difficult as n is increased. A computationally simpler algorithm is proposed next.

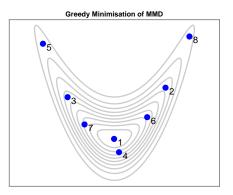


Figure: Optimal quantisation: Sequential (greedy) minimisation of MMD to select n = 8 states $\{x_i\}_{i=1}^n$ in an approximation Q_n to P. The numbers indicate the order in which the states x_i were selected.

However, the computation required to select the location x_n becomes increasingly difficult as n is increased. A computationally simpler algorithm is proposed next.

Algorithm 3 (Discrete Greedy Optimisation)

Let $(y_i)_{i\in\mathbb{N}}$ be a sample path from a Markov chain that is P-invariant. Then set

$$\mathbf{x}_n \in \underset{\mathbf{x} \in \{\mathbf{y}_1, \dots, \mathbf{y}_N\}}{\operatorname{arg\,min}} \ D_k\left(P, \frac{1}{n}\delta(\mathbf{x}) + \frac{1}{n}\sum_{i=1}^{n-1}\delta(\mathbf{x}_i)\right), \qquad n = 1, 2, \dots$$

- ▶ Includes the case where the *y_i* are independent samples from *P*. [Pre-empts application to the motivating Bayesian context.]
- ▶ For a sample path of length $N \gg 1$, the resulting sequence approximates that of continuous greedy optimisation (Algorithm 2).
- Under assumptions, that include sufficiently rapid mixing of the Markov chain

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] \leq \frac{C_1}{N} + \frac{C_2(1+\log(N))(1+\log(n))}{n}.$$

See Theorem 4 in Teymur et al. [2021].

▶ Same rate as Monte Carlo up to log terms; recall Proposition 1. [Possibly a theoretical gap here.]

Algorithm 3 (Discrete Greedy Optimisation)

Let $(y_i)_{i\in\mathbb{N}}$ be a sample path from a Markov chain that is P-invariant. Then set

$$\mathbf{x}_n \in \operatorname*{arg\,min}_{\mathbf{x} \in \{\mathbf{y}_1, \dots, \mathbf{y}_N\}} D_k \left(P, \frac{1}{n} \delta(\mathbf{x}) + \frac{1}{n} \sum_{i=1}^{n-1} \delta(\mathbf{x}_i) \right), \qquad n = 1, 2, \dots.$$

- Includes the case where the y_i are independent samples from P. [Pre-empts application to the motivating Bayesian context.]
- ▶ For a sample path of length $N \gg 1$, the resulting sequence approximates that of continuous greedy optimisation (Algorithm 2).
- Under assumptions, that include sufficiently rapid mixing of the Markov chain

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] \leq \frac{C_1}{N} + \frac{C_2(1+\log(N))(1+\log(n))}{n}.$$

See Theorem 4 in Teymur et al. [2021].

Same rate as Monte Carlo up to log terms; recall Proposition 1. [Possibly a theoretical gap here.]

Algorithm 3 (Discrete Greedy Optimisation)

Let $(y_i)_{i\in\mathbb{N}}$ be a sample path from a Markov chain that is P-invariant. Then set

$$\mathbf{x}_n \in \operatorname*{arg\,min}_{\mathbf{x} \in \{\mathbf{y}_1, \dots, \mathbf{y}_N\}} D_k \left(P, \frac{1}{n} \delta(\mathbf{x}) + \frac{1}{n} \sum_{i=1}^{n-1} \delta(\mathbf{x}_i) \right), \qquad n = 1, 2, \dots.$$

- ▶ Includes the case where the *y_i* are independent samples from *P*. [Pre-empts application to the motivating Bayesian context.]
- ▶ For a sample path of length $N \gg 1$, the resulting sequence approximates that of continuous greedy optimisation (Algorithm 2).
- Under assumptions, that include sufficiently rapid mixing of the Markov chain

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] \leq \frac{C_1}{N} + \frac{C_2(1+\log(N))(1+\log(n))}{n}.$$

See Theorem 4 in Teymur et al. [2021].

Same rate as Monte Carlo up to log terms; recall Proposition 1. [Possibly a theoretical gap here.]

Algorithm 3 (Discrete Greedy Optimisation)

Let $(y_i)_{i\in\mathbb{N}}$ be a sample path from a Markov chain that is P-invariant. Then set

$$\mathbf{x}_n \in \operatorname*{arg\,min}_{\mathbf{x} \in \{\mathbf{y}_1, \dots, \mathbf{y}_N\}} D_k \left(P, \frac{1}{n} \delta(\mathbf{x}) + \frac{1}{n} \sum_{i=1}^{n-1} \delta(\mathbf{x}_i) \right), \qquad n = 1, 2, \dots.$$

- ▶ Includes the case where the *y_i* are independent samples from *P*. [Pre-empts application to the motivating Bayesian context.]
- ▶ For a sample path of length $N \gg 1$, the resulting sequence approximates that of continuous greedy optimisation (Algorithm 2).
- Under assumptions, that include sufficiently rapid mixing of the Markov chain,

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] \leq \frac{C_1}{N} + \frac{C_2(1+\log(N))(1+\log(n))}{n}.$$

See Theorem 4 in Teymur et al. [2021].

Same rate as Monte Carlo up to log terms; recall Proposition 1. [Possibly a theoretical gap here.]

Algorithm 3 (Discrete Greedy Optimisation)

Let $(y_i)_{i\in\mathbb{N}}$ be a sample path from a Markov chain that is P-invariant. Then set

$$\mathbf{x}_n \in \operatorname*{arg\,min}_{\mathbf{x} \in \{\mathbf{y}_1, \dots, \mathbf{y}_N\}} D_k \left(P, \frac{1}{n} \delta(\mathbf{x}) + \frac{1}{n} \sum_{i=1}^{n-1} \delta(\mathbf{x}_i) \right), \qquad n = 1, 2, \dots.$$

- ▶ Includes the case where the *y_i* are independent samples from *P*. [Pre-empts application to the motivating Bayesian context.]
- ▶ For a sample path of length $N \gg 1$, the resulting sequence approximates that of continuous greedy optimisation (Algorithm 2).
- Under assumptions, that include sufficiently rapid mixing of the Markov chain,

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] \leq \frac{C_1}{N} + \frac{C_2(1+\log(N))(1+\log(n))}{n}.$$

See Theorem 4 in Teymur et al. [2021].

► Same rate as Monte Carlo up to log terms; recall Proposition 1. [Possibly a theoretical gap here.]

Algorithm 3 (Discrete Greedy Optimisation)

Let $(y_i)_{i\in\mathbb{N}}$ be a sample path from a Markov chain that is P-invariant. Then set

$$\mathbf{x}_n \in \operatorname*{arg\,min}_{\mathbf{x} \in \{\mathbf{y}_1, \dots, \mathbf{y}_N\}} D_k \left(P, \frac{1}{n} \delta(\mathbf{x}) + \frac{1}{n} \sum_{i=1}^{n-1} \delta(\mathbf{x}_i) \right), \qquad n = 1, 2, \dots.$$

- ▶ Includes the case where the *y_i* are independent samples from *P*. [Pre-empts application to the motivating Bayesian context.]
- ▶ For a sample path of length $N \gg 1$, the resulting sequence approximates that of continuous greedy optimisation (Algorithm 2).
- Under assumptions, that include sufficiently rapid mixing of the Markov chain,

$$\mathbb{E}\left[D_k(P,Q_n)^2\right] \leq \frac{C_1}{N} + \frac{C_2(1+\log(N))(1+\log(n))}{n}.$$

See Theorem 4 in Teymur et al. [2021].

▶ Same rate as Monte Carlo up to log terms; recall Proposition 1. [Possibly a theoretical gap here.]

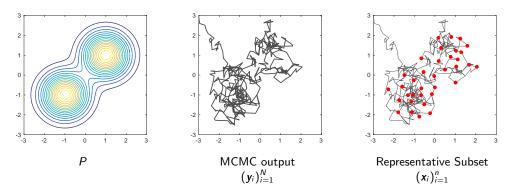


Figure: Optimal quantisation: Discrete greedy optimisation of MMD to select n=35 states $\{x_i\}_{i=1}^n$ in an approximation Q_n to P.

Q: What can be done to improve the disappointing theoretical convergence rate?

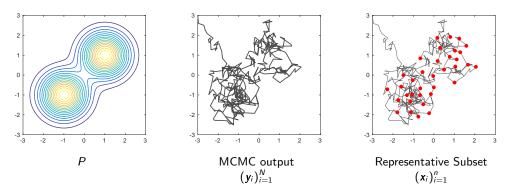


Figure: Optimal quantisation: Discrete greedy optimisation of MMD to select n=35 states $\{x_i\}_{i=1}^n$ in an approximation Q_n to P.

Q: What can be done to improve the disappointing theoretical convergence rate?

Let's return to the case of *weighted* point sets; i.e. approximations of the form $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ for some weights $w_1, \ldots, w_n \in \mathbb{R}$.

Lemma 13

Let $x_1, \ldots, x_n \in \mathcal{X}$ be distinct. The optimal weights

$$\arg\min_{\mathbf{w}\in\mathbb{R}^n}D_k\left(P,\sum_{i=1}^nw_i\delta(\mathbf{x}_i)\right)$$

are the solution of the linear system

$$Kw = z$$
 (5)

where $K_{ij} = k(\mathbf{x}_i, \mathbf{x}_j)$ and $\mathbf{z}_i = \mu_P(\mathbf{x}_i)$.

Proof

The MMD between P and $Q_n = \sum_{i=1}^n w_i \delta(x_i)$ can be expressed as

$$D_k(P,Q_n)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y)$$

Let's return to the case of weighted point sets; i.e. approximations of the form $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ for some weights $w_1, \ldots, w_n \in \mathbb{R}$.

Lemma 13

Let $x_1, \ldots, x_n \in \mathcal{X}$ be distinct. The optimal weights

$$\underset{\mathbf{w}\in\mathbb{R}^n}{\operatorname{arg\,min}}\,D_k\left(P,\sum_{i=1}^nw_i\delta(\mathbf{x}_i)\right)$$

are the solution of the linear system

$$Kw = z$$
 (5)

where $K_{ij} = k(\mathbf{x}_i, \mathbf{x}_j)$ and $\mathbf{z}_i = \mu_P(\mathbf{x}_i)$.

Proof

The MMD between P and $Q_n = \sum_{i=1}^n w_i \delta(x_i)$ can be expressed as

$$D_k(P,Q_n)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y)$$



Let's return to the case of *weighted* point sets; i.e. approximations of the form $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ for some weights $w_1, \ldots, w_n \in \mathbb{R}$.

Lemma 13

Let $x_1, \ldots, x_n \in \mathcal{X}$ be distinct. The optimal weights

$$\underset{\mathbf{w}\in\mathbb{R}^n}{\operatorname{arg\,min}}\,D_k\left(P,\sum_{i=1}^nw_i\delta(\mathbf{x}_i)\right)$$

are the solution of the linear system

$$Kw = z$$
 (5)

where $K_{ij} = k(\mathbf{x}_i, \mathbf{x}_i)$ and $\mathbf{z}_i = \mu_P(\mathbf{x}_i)$.

Proof.

The MMD between P and $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ can be expressed as

$$D_k(P,Q_n)^2 = \iint k(x,y) dP(x) dP(y) - 2 \iint k(x,y) dP(x) dQ(y) + \iint k(x,y) dQ(x) dQ(y)$$



Let's return to the case of weighted point sets; i.e. approximations of the form $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ for some weights $w_1, \ldots, w_n \in \mathbb{R}$.

Lemma 13

Let $x_1, \ldots, x_n \in \mathcal{X}$ be distinct. The optimal weights

$$\underset{\mathbf{w}\in\mathbb{R}^n}{\operatorname{arg\,min}}\,D_k\left(P,\sum_{i=1}^nw_i\delta(\mathbf{x}_i)\right)$$

are the solution of the linear system

$$Kw = z$$
 (5)

where $K_{ij} = k(\mathbf{x}_i, \mathbf{x}_j)$ and $\mathbf{z}_i = \mu_P(\mathbf{x}_i)$.

Proof.

The MMD between P and $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ can be expressed as

$$D_k(P,Q_n)^2 = C - 2\mathbf{z}^{\top}\mathbf{w} + \mathbf{w}^{\top}\mathbf{K}\mathbf{w}$$

where $C = \iint k(x, y) dP(y) dP(y)$ is independent of w. This is a non-degenerate quadratic form in w (since K is a positive definite matrix), from which the result is easily verified.

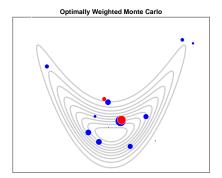


Figure: Optimally weighted Monte Carlo samples: The weights w_1, \ldots, w_n are obtained by minimising MMD in the manner of Lemma 13. Blue indicates states x_i with positive weights $w_i > 0$, while red indicates negative weights $w_i < 0$. The size of the circles is proportional to $|w_i|$.

The linear system in (5), defining the optimal weights, can be numerically ill-conditioned and the general computational overhead is $O(n^3)$. However, several "tricks" are available, such as Karvonen and Särkkä [2018], Karvonen et al. [2019].

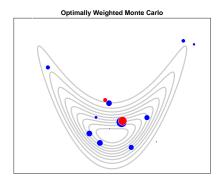


Figure: Optimally weighted Monte Carlo samples: The weights w_1, \ldots, w_n are obtained by minimising MMD in the manner of Lemma 13. Blue indicates states x_i with positive weights $w_i > 0$, while red indicates negative weights $w_i < 0$. The size of the circles is proportional to $|w_i|$.

The linear system in (5), defining the optimal weights, can be numerically ill-conditioned and the general computational overhead is $O(n^3)$. However, several "tricks" are available, such as Karvonen and Särkkä [2018], Karvonen et al. [2019].

Do optimal weights improve convergence rates?

It depends on the kernel!

Notation: The *multi-index* notation

$$\partial^{\alpha} f: x \mapsto \frac{\partial^{\alpha} f(x)}{\partial x^{\alpha}} = \frac{\partial^{|\alpha|}}{\partial x_1^{\alpha_1} \dots \partial x_d^{\alpha_d}} f(x), \qquad \alpha \in \mathbb{N}_0^d$$

will be used, and we let $|\alpha| = \alpha_1 + \cdots + \alpha_d$.

Definition 14

For s > d/2 and (sufficiently regular) $\mathcal{X} \subset \mathbb{R}^d$, the (order s) Sobolev space $H^s(\mathcal{X})$ is defined to be the set of functions $f: \mathcal{X} \to \mathbb{R}$ whose mixed partial derivatives $\partial^{\alpha} f$, $|\alpha| \leq s$, exist in $L^2(\mathcal{X})$. This becomes a Hilbert space with inner product

$$\langle f, g \rangle_{H^s(\mathcal{X})} = \sum_{|\alpha| \le s} \int \frac{\partial^{\alpha} f(x)}{\partial x^{\alpha}} \frac{\partial^{\alpha} g(x)}{\partial x^{\alpha}} dx$$

Do optimal weights improve convergence rates?

It depends on the kernel!

Notation: The *multi-index* notation

$$\partial^{\alpha} f: x \mapsto \frac{\partial^{\alpha} f(x)}{\partial x^{\alpha}} = \frac{\partial^{|\alpha|}}{\partial x_1^{\alpha_1} \dots \partial x_d^{\alpha_d}} f(x), \qquad \alpha \in \mathbb{N}_0^d$$

will be used, and we let $|\alpha| = \alpha_1 + \cdots + \alpha_d$.

Definition 14

For s > d/2 and (sufficiently regular) $\mathcal{X} \subset \mathbb{R}^d$, the (order s) Sobolev space $H^s(\mathcal{X})$ is defined to be the set of functions $f: \mathcal{X} \to \mathbb{R}$ whose mixed partial derivatives $\partial^{\alpha} f$, $|\alpha| \leq s$, exist in $L^2(\mathcal{X})$. This becomes a Hilbert space with inner product

$$\langle f, g \rangle_{H^s(\mathcal{X})} = \sum_{|\alpha| \le s} \int \frac{\partial^{\alpha} f(x)}{\partial x^{\alpha}} \frac{\partial^{\alpha} g(x)}{\partial x^{\alpha}} dx$$

It depends on the kernel!

Notation: The multi-index notation

$$\partial^{\alpha} f: \mathbf{x} \mapsto \frac{\partial^{\alpha} f(\mathbf{x})}{\partial \mathbf{x}^{\alpha}} = \frac{\partial^{|\alpha|}}{\partial x_{1}^{\alpha_{1}} \dots \partial x_{d}^{\alpha_{d}}} f(\mathbf{x}), \qquad \alpha \in \mathbb{N}_{0}^{d}$$

will be used, and we let $|\alpha| = \alpha_1 + \cdots + \alpha_d$.

Definition 14

For s>d/2 and (sufficiently regular) $\mathcal{X}\subset\mathbb{R}^d$, the (order s) Sobolev space $H^s(\mathcal{X})$ is defined to be the set of functions $f:\mathcal{X}\to\mathbb{R}$ whose mixed partial derivatives $\partial^{\alpha}f$, $|\alpha|\leq s$, exist in $L^2(\mathcal{X})$. This becomes a Hilbert space with inner product

$$\langle f, g \rangle_{H^s(\mathcal{X})} = \sum_{|\alpha| \le s} \int \frac{\partial^{\alpha} f(x)}{\partial x^{\alpha}} \frac{\partial^{\alpha} g(x)}{\partial x^{\alpha}} dx$$

It depends on the kernel!

Notation: The multi-index notation

$$\partial^{\alpha} f: \mathbf{x} \mapsto \frac{\partial^{\alpha} f(\mathbf{x})}{\partial \mathbf{x}^{\alpha}} = \frac{\partial^{|\alpha|}}{\partial x_{1}^{\alpha_{1}} \dots \partial x_{d}^{\alpha_{d}}} f(\mathbf{x}), \qquad \alpha \in \mathbb{N}_{0}^{d}$$

will be used, and we let $|\alpha| = \alpha_1 + \cdots + \alpha_d$.

Definition 14

For s>d/2 and (sufficiently regular) $\mathcal{X}\subset\mathbb{R}^d$, the (order s) Sobolev space $H^s(\mathcal{X})$ is defined to be the set of functions $f:\mathcal{X}\to\mathbb{R}$ whose mixed partial derivatives $\partial^{\alpha}f$, $|\alpha|\leq s$, exist in $L^2(\mathcal{X})$. This becomes a Hilbert space with inner product

$$\langle f, g \rangle_{H^s(\mathcal{X})} = \sum_{|\alpha| \leq s} \int \frac{\partial^{\alpha} f(x)}{\partial x^{\alpha}} \frac{\partial^{\alpha} g(x)}{\partial x^{\alpha}} dx.$$

It depends on the kernel!

Notation: The multi-index notation

$$\partial^{\alpha} f: \mathbf{x} \mapsto \frac{\partial^{\alpha} f(\mathbf{x})}{\partial \mathbf{x}^{\alpha}} = \frac{\partial^{|\alpha|}}{\partial x_{1}^{\alpha_{1}} \dots \partial x_{d}^{\alpha_{d}}} f(\mathbf{x}), \qquad \alpha \in \mathbb{N}_{0}^{d}$$

will be used, and we let $|\alpha| = \alpha_1 + \cdots + \alpha_d$.

Definition 14

For s>d/2 and (sufficiently regular) $\mathcal{X}\subset\mathbb{R}^d$, the (order s) Sobolev space $H^s(\mathcal{X})$ is defined to be the set of functions $f:\mathcal{X}\to\mathbb{R}$ whose mixed partial derivatives $\partial^{\alpha}f$, $|\alpha|\leq s$, exist in $L^2(\mathcal{X})$. This becomes a Hilbert space with inner product

$$\langle f, g \rangle_{H^s(\mathcal{X})} = \sum_{|\alpha| \le s} \int \frac{\partial^{\alpha} f(x)}{\partial x^{\alpha}} \frac{\partial^{\alpha} g(x)}{\partial x^{\alpha}} dx.$$

It depends on the kernel!

Notation: The *multi-index* notation

$$\partial^{\alpha} f: \mathbf{x} \mapsto \frac{\partial^{\alpha} f(\mathbf{x})}{\partial \mathbf{x}^{\alpha}} = \frac{\partial^{|\alpha|}}{\partial x_{1}^{\alpha_{1}} \dots \partial x_{d}^{\alpha_{d}}} f(\mathbf{x}), \qquad \alpha \in \mathbb{N}_{0}^{d}$$

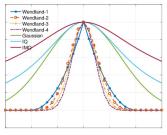
will be used, and we let $|\alpha| = \alpha_1 + \cdots + \alpha_d$.

Example 14

Let z_+^m denote $\max(0,z)^m$ in shorthand. Then examples of Sobolev kernels on (sufficiently regular) $\mathcal{X} \subseteq \mathbb{R}$ include the following, due to Wendland [1998]:

$k(x,y)$ $(r= x-y , x,y \in \mathbb{R})$	order
$(1-r)_{+}$	s=1
$(1-r)^3_+(3r+1)$	s=2
$(1-r)^{5}_{+}(8r^{2}+5r+1)$	<i>s</i> = 3

These kernels are convenient for numerical reasons, due to their compact support, which renders K a sparse matrix.



Optimal weights can accelerate convergence:

Theorem 15 (e.g. Ehler et al. [2019])

Let $x_1, \ldots, x_n \sim P$ be independent and let $\mathbf{w} = \mathbf{w}(x_1, \ldots, x_n)$ denote optimal weights in the sense of Lemma 13. Let $k(\mathbf{x}, \mathbf{y})$ be an (order s) Sobolev kernel. Then, under regularity conditions on the domain \mathcal{X} , which is of dimension d, and the distribution P, there exists a constant $0 < C < \infty$ such that

$$\mathbb{E}\left[D_k\left(P,\sum_{i=1}^n w_i\delta(\mathbf{x}_i)\right)\right] \leq C\left(\frac{\log(n)}{n}\right)^{s/d}.$$

For s = d/2 we recover the same rate as Proposition 1 for un-weighted Monte Carlo (up to log factors), while for s > d/2 we obtain faster convergence in MMD.

So surely it is a good idea to employ optimal weights? Not necessarily - the gain in rate has to out-weight the higher computational cost of $O(n^3)$. Superior error-per-cost only when s > 3d/2

Numerical ill-conditioning at large n also requires careful treatment. The use of points from optimal quantisation can be helpful.



Optimal weights can accelerate convergence:

Theorem 15 (e.g. Ehler et al. [2019])

Let $x_1, \ldots, x_n \sim P$ be independent and let $\mathbf{w} = \mathbf{w}(x_1, \ldots, x_n)$ denote optimal weights in the sense of Lemma 13. Let $k(\mathbf{x}, \mathbf{y})$ be an (order s) Sobolev kernel. Then, under regularity conditions on the domain \mathcal{X} , which is of dimension d, and the distribution P, there exists a constant $0 < C < \infty$ such that

$$\mathbb{E}\left[D_k\left(P,\sum_{i=1}^n w_i\delta(\mathbf{x}_i)\right)\right] \leq C\left(\frac{\log(n)}{n}\right)^{s/d}.$$

For s=d/2 we recover the same rate as Proposition 1 for un-weighted Monte Carlo (up to log factors), while for s>d/2 we obtain faster convergence in MMD.

So surely it is a good idea to employ optimal weights? Not necessarily - the gain in rate has to out-weight the higher computational cost of $O(n^3)$. Superior error-per-cost only when s > 3d/2

Numerical ill-conditioning at large n also requires careful treatment. The use of points from optimal quantisation can be helpful.



Optimal weights <u>can</u> accelerate convergence:

Theorem 15 (e.g. Ehler et al. [2019])

Let $x_1, \ldots, x_n \sim P$ be independent and let $\mathbf{w} = \mathbf{w}(x_1, \ldots, x_n)$ denote optimal weights in the sense of Lemma 13. Let $k(\mathbf{x}, \mathbf{y})$ be an (order s) Sobolev kernel. Then, under regularity conditions on the domain \mathcal{X} , which is of dimension d, and the distribution P, there exists a constant $0 < C < \infty$ such that

$$\mathbb{E}\left[D_k\left(P,\sum_{i=1}^n w_i\delta(\mathbf{x}_i)\right)\right] \leq C\left(\frac{\log(n)}{n}\right)^{s/d}.$$

For s=d/2 we recover the same rate as Proposition 1 for un-weighted Monte Carlo (up to log factors), while for s>d/2 we obtain faster convergence in MMD.

So surely it is a good idea to employ optimal weights? Not necessarily - the gain in rate has to out-weight the higher computational cost of $O(n^3)$. Superior error-per-cost only when s > 3d/2

Numerical ill-conditioning at large n also requires careful treatment. The use of points from optimal quantisation can be helpful.



Optimal weights <u>can</u> accelerate convergence:

Theorem 15 (e.g. Ehler et al. [2019])

Let $x_1, \ldots, x_n \sim P$ be independent and let $\mathbf{w} = \mathbf{w}(x_1, \ldots, x_n)$ denote optimal weights in the sense of Lemma 13. Let $k(\mathbf{x}, \mathbf{y})$ be an (order s) Sobolev kernel. Then, under regularity conditions on the domain \mathcal{X} , which is of dimension d, and the distribution P, there exists a constant $0 < C < \infty$ such that

$$\mathbb{E}\left[D_k\left(P,\sum_{i=1}^n w_i\delta(\mathbf{x}_i)\right)\right] \leq C\left(\frac{\log(n)}{n}\right)^{s/d}.$$

For s = d/2 we recover the same rate as Proposition 1 for un-weighted Monte Carlo (up to log factors), while for s > d/2 we obtain faster convergence in MMD.

So surely it is a good idea to employ optimal weights? Not necessarily - the gain in rate has to out-weight the higher computational cost of $O(n^3)$. Superior error-per-cost only when s > 3d/2.

Numerical ill-conditioning at large n also requires careful treatment. The use of points from optimal quantisation can be helpful.



Optimal weights <u>can</u> accelerate convergence:

Theorem 15 (e.g. Ehler et al. [2019])

Let $x_1, \ldots, x_n \sim P$ be independent and let $\mathbf{w} = \mathbf{w}(x_1, \ldots, x_n)$ denote optimal weights in the sense of Lemma 13. Let $k(\mathbf{x}, \mathbf{y})$ be an (order s) Sobolev kernel. Then, under regularity conditions on the domain \mathcal{X} , which is of dimension d, and the distribution P, there exists a constant $0 < C < \infty$ such that

$$\mathbb{E}\left[D_k\left(P,\sum_{i=1}^n w_i\delta(\mathbf{x}_i)\right)\right] \leq C\left(\frac{\log(n)}{n}\right)^{s/d}.$$

For s=d/2 we recover the same rate as Proposition 1 for un-weighted Monte Carlo (up to log factors), while for s>d/2 we obtain faster convergence in MMD.

So surely it is a good idea to employ optimal weights? Not necessarily - the gain in rate has to out-weight the higher computational cost of $O(n^3)$. Superior error-per-cost only when s > 3d/2.

Numerical ill-conditioning at large n also requires careful treatment. The use of points from optimal quantisation can be helpful.



Optimal weights <u>can</u> accelerate convergence:

Theorem 15 (e.g. Ehler et al. [2019])

Let $x_1, \ldots, x_n \sim P$ be independent and let $\mathbf{w} = \mathbf{w}(x_1, \ldots, x_n)$ denote optimal weights in the sense of Lemma 13. Let $k(\mathbf{x}, \mathbf{y})$ be an (order s) Sobolev kernel. Then, under regularity conditions on the domain \mathcal{X} , which is of dimension d, and the distribution P, there exists a constant $0 < C < \infty$ such that

$$\mathbb{E}\left[D_k\left(P,\sum_{i=1}^n w_i\delta(\mathbf{x}_i)\right)\right] \leq C\left(\frac{\log(n)}{n}\right)^{s/d}.$$

For s=d/2 we recover the same rate as Proposition 1 for un-weighted Monte Carlo (up to log factors), while for s>d/2 we obtain faster convergence in MMD.

So surely it is a good idea to employ optimal weights? Not necessarily - the gain in rate has to out-weight the higher computational cost of $O(n^3)$. Superior error-per-cost only when s > 3d/2.

Numerical ill-conditioning at large n also requires careful treatment. The use of points from optimal quantisation can be helpful.



Recall that we aim to perform optimal quantisation of a distribution P that admits a PDF p(x) on $x \in \mathbb{R}^d$, such that

$$p(x)=\frac{\tilde{p}(x)}{Z},$$

where \tilde{p} can be exactly evaluated but Z, and hence p(x), cannot easily be evaluated or even approximated.

This setting is typical in applications of Bayesian inference, where we have

$$p(x) = \frac{\pi(x)\mathcal{L}(x)}{Z}$$

where $\pi(x)$ is a prior PDF, $\mathcal{L}(x)$ is a likelihood, and the implicitly defined normalisation constant Z is the marginal likelihood.

Several methods have been developed in the statistics, applied probability, physics and machine learning literatures to approximate distributions P with these characteristics, including Markov chain Monte Carlo (MCMC), sequential Monte Carlo (SMC), and Variational inference. These techniques do not typically attempt Optimal quantisation, since even the basic quantisation task can be difficult.

Recall that we aim to perform optimal quantisation of a distribution P that admits a PDF p(x) on $x \in \mathbb{R}^d$, such that

$$p(x)=\frac{\tilde{p}(x)}{Z},$$

where \tilde{p} can be exactly evaluated but Z, and hence p(x), cannot easily be evaluated or even approximated.

This setting is typical in applications of Bayesian inference, where we have

$$p(x) = \frac{\pi(x)\mathcal{L}(x)}{Z}$$

where $\pi(x)$ is a prior PDF, $\mathcal{L}(x)$ is a likelihood, and the implicitly defined normalisation constant Z is the marginal likelihood.

Several methods have been developed in the statistics, applied probability, physics and machine learning literatures to approximate distributions P with these characteristics, including MCMC, SMC, and variational inference. These techniques do not typically attempt optimal quantisation, since even the basic quantisation task can be difficult.

Recall that we aim to perform optimal quantisation of a distribution P that admits a PDF p(x) on $x \in \mathbb{R}^d$, such that

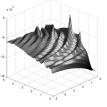
$$p(x)=\frac{\tilde{p}(x)}{Z},$$

where \tilde{p} can be exactly evaluated but Z, and hence p(x), cannot easily be evaluated or even approximated.

The integral

$$Z = \int \pi(\mathbf{x}) \mathcal{L}(\mathbf{x}) \mathrm{d}\mathbf{x}$$

is often extremely challenging to evaluate due to localised regions in which $\ensuremath{\mathcal{L}}$ takes very large values.



Several methods have been developed in the statistics, applied probability, physics and machine learning literatures to approximate distributions *P* with these characteristics, including *MCMC*, *SMC*, and *variational inference*. These techniques do not typically attempt *optimal* quantisation, since even the basic quantisation task can be difficult.

Recall that we aim to perform optimal quantisation of a distribution P that admits a PDF p(x) on $x \in \mathbb{R}^d$, such that

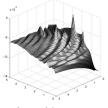
$$p(x)=\frac{\tilde{p}(x)}{Z},$$

where \tilde{p} can be exactly evaluated but Z, and hence p(x), cannot easily be evaluated or even approximated.

The integral

$$Z = \int \pi(\mathbf{x}) \mathcal{L}(\mathbf{x}) \mathrm{d}\mathbf{x}$$

is often extremely challenging to evaluate due to localised regions in which $\ensuremath{\mathcal{L}}$ takes very large values.



Several methods have been developed in the statistics, applied probability, physics and machine learning literatures to approximate distributions *P* with these characteristics, including *MCMC*, *SMC*, and *variational inference*. These techniques do not typically attempt *optimal* quantisation, since even the basic quantisation task can be difficult.

How can our algorithms be applied?

The apparent difficulty is that we cannot compute integrals with respect to P, such as $\int k(\cdot, x) dP(x)$, which are required for computation of MMD. A hint at a possible solution is provided by the following result:

Lemma 1

Suppose $k_P:\mathcal{X} imes\mathcal{X} o\mathbb{R}$ is a kernel with $\int k_P(\cdot,x)\mathrm{d}P=0$ for all $x\in\mathcal{X}$. Ther

$$D_{k_P}(Q) = D_{k_P}(P,Q) = \sup_{\|f\|_{\mathcal{H}(k_P)} \le 1} \left| \int f dQ \right|$$

Proof

For all $f \in \mathcal{H}(k_P)$ it holds that $\int f \, \mathrm{d}P = 0$, whence the result. Indeed, from the reproducing property, and using Lemma 6 with Standing Assumption 1 to interchange integral with inner product, $\int f \, \mathrm{d}P = \int \langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k_P)} \, \mathrm{d}P(\mathbf{x}) = \langle f, f \rangle_{\mathcal{H}(k_P)} = \langle f, f \rangle_{\mathcal{H}(k_P)} = 0$.

How can our algorithms be applied?

The apparent difficulty is that we cannot compute integrals with respect to P, such as $\int k(\cdot, x) dP(x)$, which are required for computation of MMD. A hint at a possible solution is provided by the following result:

Lemma 16

Suppose $k_P: \mathcal{X} \times \mathcal{X} \to \mathbb{R}$ is a kernel with $\int k_P(\cdot, \mathbf{x}) dP = 0$ for all $\mathbf{x} \in \mathcal{X}$. Then

$$D_{k_P}(Q) = D_{k_P}(P,Q) = \sup_{\|f\|_{\mathcal{H}(k_P)} \le 1} \left| \int f \mathrm{d}Q \right|.$$

Proof

For all $f \in \mathcal{H}(k_P)$ it holds that $\int f \, \mathrm{d}P = 0$, whence the result. Indeed, from the reproducing property, and using Lemma 6 with Standing Assumption 1 to interchange integral with inner product, $\int f \, \mathrm{d}P = \int \langle f, k(\cdot, x) \rangle_{\mathcal{H}(k_P)} \, \mathrm{d}P(x) = \langle f, \int k(\cdot, x) \, \mathrm{d}P(x) \rangle_{\mathcal{H}(k_P)} = 0$.

How can our algorithms be applied?

The apparent difficulty is that we cannot compute integrals with respect to P, such as $\int k(\cdot, x) dP(x)$, which are required for computation of MMD. A hint at a possible solution is provided by the following result:

Lemma 16

Suppose $k_P: \mathcal{X} \times \mathcal{X} \to \mathbb{R}$ is a kernel with $\int k_P(\cdot, \mathbf{x}) dP = 0$ for all $\mathbf{x} \in \mathcal{X}$. Then

$$D_{k_P}(Q) = D_{k_P}(P,Q) = \sup_{\|f\|_{\mathcal{H}(k_P)} \le 1} \left| \int f dQ \right|$$

Proof

For all $f \in \mathcal{H}(k_P)$ it holds that $\int f dP = 0$, whence the result. Indeed, from the reproducing property, and using Lemma 6 with Standing Assumption 1 to interchange integral with inner product, $\int f dP = \int \langle f, k(\cdot, \mathbf{x}) \rangle_{2M_P} dP(\mathbf{x}) = \langle f, f, k(\cdot, \mathbf{x}) \rangle_{2M_P} dP(\mathbf{x}) = \langle f, 0 \rangle_{2M_P} = 0$

How can our algorithms be applied?

The apparent difficulty is that we cannot compute integrals with respect to P, such as $\int k(\cdot, x) dP(x)$, which are required for computation of MMD. A hint at a possible solution is provided by the following result:

Lemma 16

Suppose $k_P: \mathcal{X} \times \mathcal{X} \to \mathbb{R}$ is a kernel with $\int k_P(\cdot, \mathbf{x}) \mathrm{d}P = 0$ for all $\mathbf{x} \in \mathcal{X}$. Then

$$D_{k_P}(Q) = D_{k_P}(P,Q) = \sup_{\|f\|_{\mathcal{H}(k_P)} \le 1} \left| \int f dQ \right|.$$

Proof

For all $f \in \mathcal{H}(k_P)$ it holds that $\int f dP = 0$, whence the result. Indeed, from the reproducing property, and using Lemma 6 with Standing Assumption 1 to interchange integral with inner product, $\int f dP = \int \langle f | k(\cdot, \mathbf{x}) \rangle_{AUV} dP(\mathbf{x}) = \langle f | f | k(\cdot, \mathbf{x}) dP(\mathbf{x}) \rangle_{AUV} = \langle f | f \rangle_{AUV} = 0$

How can our algorithms be applied?

The apparent difficulty is that we cannot compute integrals with respect to P, such as $\int k(\cdot, x) dP(x)$, which are required for computation of MMD. A hint at a possible solution is provided by the following result:

Lemma 16

Suppose $k_P: \mathcal{X} \times \mathcal{X} \to \mathbb{R}$ is a kernel with $\int k_P(\cdot, \mathbf{x}) \mathrm{d}P = 0$ for all $\mathbf{x} \in \mathcal{X}$. Then

$$D_{k_P}(Q) = D_{k_P}(P,Q) = \sup_{\|f\|_{\mathcal{H}(k_P)} \le 1} \left| \int f dQ \right|.$$

Proof.

For all $f \in \mathcal{H}(k_P)$ it holds that $\int f \mathrm{d}P = 0$, whence the result. Indeed, from the reproducing property, and using Lemma 6 with Standing Assumption 1 to interchange integral with inner product, $\int f \mathrm{d}P = \int \langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k_P)} \mathrm{d}P(\mathbf{x}) = \left\langle f, \int k(\cdot, \mathbf{x}) \mathrm{d}P(\mathbf{x}) \right\rangle_{\mathcal{H}(k_P)} = \langle f, 0 \rangle_{\mathcal{H}(k_P)} = 0.$



How can our algorithms be applied?

The apparent difficulty is that we cannot compute integrals with respect to P, such as $\int k(\cdot, x) dP(x)$, which are required for computation of MMD. A hint at a possible solution is provided by the following result:

Lemma 16

Suppose $k_P: \mathcal{X} \times \mathcal{X} \to \mathbb{R}$ is a kernel with $\int k_P(\cdot, \mathbf{x}) dP = 0$ for all $\mathbf{x} \in \mathcal{X}$. Then

$$D_{k_P}(Q) = D_{k_P}(P,Q) = \sup_{\|f\|_{\mathcal{H}(k_P)} \le 1} \left| \int f dQ \right|.$$

Proof.

For all $f \in \mathcal{H}(k_P)$ it holds that $\int f \, \mathrm{d}P = 0$, whence the result. Indeed, from the reproducing property, and using Lemma 6 with Standing Assumption 1 to interchange integral with inner product, $\int f \, \mathrm{d}P = \int \langle f, k(\cdot, \mathbf{x}) \rangle_{\mathcal{H}(k_P)} \, \mathrm{d}P(\mathbf{x}) = \langle f, \int k(\cdot, \mathbf{x}) \, \mathrm{d}P(\mathbf{x}) \rangle_{\mathcal{H}(k_P)} = \langle f, 0 \rangle_{\mathcal{H}(k_P)} = 0.$

As an example, consider a bounded linear operator

$$(\mathcal{A}_P g)(x) = g(x) - \int g dP$$

acting on elements of an RKHS $\mathcal{H}(k)$.

If we apply A_P to both arguments of the kernel k, we obtain a Stein kernel

$$k_{P}(x,y) = \mathcal{A}_{P}^{y} \mathcal{A}_{P}^{x} k(x,y)$$

$$= k(x,y) - \int k(x,y) dP(x) - \int k(x,y) dP(y) + \iint k(x,y) dP(x) dP(y). \tag{6}$$

Indeed, $\int k_P(\cdot, x) dP(x) = \int \mathcal{A}_P^y \mathcal{A}_P^x k(x, y) dP(x) = \mathcal{A}_P^y \int \mathcal{A}_P^x k(x, y) dP(x) = \mathcal{A}_P^y 0 = 0$, where interchange of \mathcal{A}_P^y and the integral is justified by noting that \mathcal{A}_P^y is a bounded linear operator and following similar reasoning to Lemma 6.

Unfortunately, the Stein kernel in (6) is not useful because it still involves the problematic integral $\int k(\cdot, x) dP(x)$. A more useful construction of a Stein kernel is needed.

As an example, consider a bounded linear operator

$$(\mathcal{A}_P g)(x) = g(x) - \int g dP$$

acting on elements of an RKHS $\mathcal{H}(k)$.

If we apply A_P to both arguments of the kernel k, we obtain a Stein kernel

$$k_{P}(x,y) = A_{P}^{y} A_{P}^{x} k(x,y)$$

$$= k(x,y) - \int k(x,y) dP(x) - \int k(x,y) dP(y) + \iint k(x,y) dP(x) dP(y).$$
(6)

Indeed, $\int k_P(\cdot, \mathbf{x}) dP(\mathbf{x}) = \int \mathcal{A}_P^{\mathbf{y}} \mathcal{A}_P^{\mathbf{x}} k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) = \mathcal{A}_P^{\mathbf{y}} \int \mathcal{A}_P^{\mathbf{x}} k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) = \mathcal{A}_P^{\mathbf{y}} 0 = 0$, where interchange of $\mathcal{A}_P^{\mathbf{y}}$ and the integral is justified by noting that $\mathcal{A}_P^{\mathbf{y}}$ is a bounded linear operator and following similar reasoning to Lemma 6.

Unfortunately, the Stein kernel in (6) is not useful because it still involves the problematic integral $\int k(\cdot, x) dP(x)$. A more useful construction of a Stein kernel is needed.

As an example, consider a bounded linear operator

$$(\mathcal{A}_P g)(x) = g(x) - \int g dP$$

acting on elements of an RKHS $\mathcal{H}(k)$.

If we apply A_P to both arguments of the kernel k, we obtain a Stein kernel

$$k_{P}(\mathbf{x}, \mathbf{y}) = \mathcal{A}_{P}^{\mathbf{y}} \mathcal{A}_{P}^{\mathbf{x}} k(\mathbf{x}, \mathbf{y})$$

$$= k(\mathbf{x}, \mathbf{y}) - \int k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) - \int k(\mathbf{x}, \mathbf{y}) dP(\mathbf{y}) + \iint k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) dP(\mathbf{y}). \tag{6}$$

Indeed, $\int k_P(\cdot, \mathbf{x}) dP(\mathbf{x}) = \int \mathcal{A}_P^{\mathbf{y}} \mathcal{A}_P^{\mathbf{x}} k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) = \mathcal{A}_P^{\mathbf{y}} \int \mathcal{A}_P^{\mathbf{x}} k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) = \mathcal{A}_P^{\mathbf{y}} 0 = 0$, where interchange of $\mathcal{A}_P^{\mathbf{y}}$ and the integral is justified by noting that $\mathcal{A}_P^{\mathbf{y}}$ is a bounded linear operator and following similar reasoning to Lemma 6.

Unfortunately, the Stein kernel in (6) is not useful because it still involves the problematic integral $\int k(\cdot, x) dP(x)$. A more useful construction of a Stein kernel is needed.

Stein Discrepancy

As an example, consider a bounded linear operator

$$(\mathcal{A}_P g)(x) = g(x) - \int g dP$$

acting on elements of an RKHS $\mathcal{H}(k)$.

If we apply A_P to both arguments of the kernel k, we obtain a Stein kernel

$$k_{P}(\mathbf{x}, \mathbf{y}) = \mathcal{A}_{P}^{\mathbf{y}} \mathcal{A}_{P}^{\mathbf{x}} k(\mathbf{x}, \mathbf{y})$$

$$= k(\mathbf{x}, \mathbf{y}) - \int k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) - \int k(\mathbf{x}, \mathbf{y}) dP(\mathbf{y}) + \iint k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) dP(\mathbf{y}). \tag{6}$$

Indeed, $\int k_P(\cdot, x) dP(x) = \int \mathcal{A}_P^y \mathcal{A}_P^x k(x, y) dP(x) = \mathcal{A}_P^y \int \mathcal{A}_P^x k(x, y) dP(x) = \mathcal{A}_P^y 0 = 0$, where interchange of \mathcal{A}_P^y and the integral is justified by noting that \mathcal{A}_P^y is a bounded linear operator and following similar reasoning to Lemma 6.

Unfortunately, the Stein kernel in (6) is not useful because it still involves the problematic integral $\int k(\cdot, x) dP(x)$. A more useful construction of a Stein kernel is needed.

Stein Discrepancy

As an example, consider a bounded linear operator

$$(\mathcal{A}_P g)(x) = g(x) - \int g dP$$

acting on elements of an RKHS $\mathcal{H}(k)$.

If we apply A_P to both arguments of the kernel k, we obtain a Stein kernel

$$k_{P}(\mathbf{x}, \mathbf{y}) = \mathcal{A}_{P}^{\mathbf{y}} \mathcal{A}_{P}^{\mathbf{x}} k(\mathbf{x}, \mathbf{y})$$

$$= k(\mathbf{x}, \mathbf{y}) - \int k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) - \int k(\mathbf{x}, \mathbf{y}) dP(\mathbf{y}) + \iint k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) dP(\mathbf{y}). \tag{6}$$

Indeed, $\int k_P(\cdot, \mathbf{x}) dP(\mathbf{x}) = \int \mathcal{A}_P^{\mathbf{y}} \mathcal{A}_P^{\mathbf{x}} k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) = \mathcal{A}_P^{\mathbf{y}} \int \mathcal{A}_P^{\mathbf{x}} k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) = \mathcal{A}_P^{\mathbf{y}} 0 = 0$, where interchange of $\mathcal{A}_P^{\mathbf{y}}$ and the integral is justified by noting that $\mathcal{A}_P^{\mathbf{y}}$ is a bounded linear operator and following similar reasoning to Lemma 6.

Unfortunately, the Stein kernel in (6) is not useful because it still involves the problematic integral $\int k(\cdot, x) dP(x)$. A more useful construction of a Stein kernel is needed.

Stein Discrepancy

As an example, consider a bounded linear operator

$$(\mathcal{A}_P g)(x) = g(x) - \int g dP$$

acting on elements of an RKHS $\mathcal{H}(k)$.

If we apply A_P to both arguments of the kernel k, we obtain a Stein kernel

$$k_{P}(\mathbf{x}, \mathbf{y}) = \mathcal{A}_{P}^{\mathbf{y}} \mathcal{A}_{P}^{\mathbf{x}} k(\mathbf{x}, \mathbf{y})$$

$$= k(\mathbf{x}, \mathbf{y}) - \int k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) - \int k(\mathbf{x}, \mathbf{y}) dP(\mathbf{y}) + \iint k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) dP(\mathbf{y}). \tag{6}$$

Indeed, $\int k_P(\cdot, \mathbf{x}) dP(\mathbf{x}) = \int \mathcal{A}_P^{\mathbf{y}} \mathcal{A}_P^{\mathbf{x}} k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) = \mathcal{A}_P^{\mathbf{y}} \int \mathcal{A}_P^{\mathbf{x}} k(\mathbf{x}, \mathbf{y}) dP(\mathbf{x}) = \mathcal{A}_P^{\mathbf{y}} 0 = 0$, where interchange of $\mathcal{A}_P^{\mathbf{y}}$ and the integral is justified by noting that $\mathcal{A}_P^{\mathbf{y}}$ is a bounded linear operator and following similar reasoning to Lemma 6.

Unfortunately, the Stein kernel in (6) is not useful because it still involves the problematic integral $\int k(\cdot, x) dP(x)$. A more useful construction of a Stein kernel is needed.

Notation: Let $\nabla f = (\partial_{x_1} f, \dots, \partial_{x_d} f)^{\top}$ for differentiable functions $f : \mathbb{R}^d \to \mathbb{R}$.

Standing Assumption 2

The distribution P admits a positive and differentiable PDF such that $x \mapsto (\nabla \log p)(x)$ is Lipschitz.

Definition 17 (Canonical Stein operator)

For a distribution P admitting a positive and differentiable density p on \mathbb{R}^d , we define the *canonical Stein operator*

$$(\mathcal{A}_{P}g)(x) = (\nabla \cdot g)(x) + g(x) \cdot (\nabla \log p)(x)$$

acting on differentiable vector field $g : \mathbb{R}^d \to \mathbb{R}^d$, where $\mathbf{x} \in \mathbb{R}^d$.

The canonical Stein operator was introduced (for Gaussian P) in Stein [1972]. Importantly, observe that

$$(\nabla \log p)(x) = \frac{(\nabla p)(x)}{p(x)} = \frac{\frac{1}{Z}(\nabla \tilde{p})(x)}{\frac{1}{Z}\tilde{p}(x)} = \frac{(\nabla \tilde{p})(x)}{\tilde{p}(x)} = (\nabla \log \tilde{p})(x)$$

which can be computed without knowledge of p or Z, provided \tilde{p} and $\nabla \tilde{p}$ can be evaluated

Notation: Let $\nabla f = (\partial_{x_1} f, \dots, \partial_{x_d} f)^{\top}$ for differentiable functions $f : \mathbb{R}^d \to \mathbb{R}$.

Standing Assumption 2

The distribution P admits a positive and differentiable PDF such that $x \mapsto (\nabla \log p)(x)$ is Lipschitz.

Definition 17 (Canonical Stein operator)

For a distribution P admitting a positive and differentiable density p on \mathbb{R}^d , we define the *canonical Stein operator*

$$(A_P g)(x) = (\nabla \cdot g)(x) + g(x) \cdot (\nabla \log p)(x)$$

acting on differentiable vector field $g : \mathbb{R}^d \to \mathbb{R}^d$, where $\mathbf{x} \in \mathbb{R}^d$.

The canonical Stein operator was introduced (for Gaussian P) in Stein [1972]. Importantly, observe that

$$(\nabla \log p)(x) = \frac{(\nabla p)(x)}{p(x)} = \frac{\frac{1}{Z}(\nabla \tilde{p})(x)}{\frac{1}{Z}\tilde{p}(x)} = \frac{(\nabla \tilde{p})(x)}{\tilde{p}(x)} = (\nabla \log \tilde{p})(x)$$

which can be computed without knowledge of p or Z, provided \tilde{p} and $\nabla \tilde{p}$ can be evaluated

Notation: Let $\nabla f = (\partial_{x_1} f, \dots, \partial_{x_d} f)^{\top}$ for differentiable functions $f : \mathbb{R}^d \to \mathbb{R}$.

Standing Assumption 2

The distribution P admits a positive and differentiable PDF such that $x \mapsto (\nabla \log p)(x)$ is Lipschitz.

Definition 17 (Canonical Stein operator)

For a distribution P admitting a positive and differentiable density p on \mathbb{R}^d , we define the *canonical Stein operator*

$$(\mathcal{A}_P g)(x) = (\nabla \cdot g)(x) + g(x) \cdot (\nabla \log p)(x)$$

acting on differentiable vector field $g : \mathbb{R}^d \to \mathbb{R}^d$, where $\mathbf{x} \in \mathbb{R}^d$.

The canonical Stein operator was introduced (for Gaussian P) in Stein [1972]. Importantly, observe that

$$(\nabla \log p)(x) = \frac{(\nabla p)(x)}{p(x)} = \frac{\frac{1}{Z}(\nabla \tilde{p})(x)}{\frac{1}{Z}\tilde{p}(x)} = \frac{(\nabla \tilde{p})(x)}{\tilde{p}(x)} = (\nabla \log \tilde{p})(x)$$

which can be computed without knowledge of p or Z, provided \tilde{p} and $\nabla \tilde{p}$ can be evaluated.



Notation: Let $\nabla f = (\partial_{x_1} f, \dots, \partial_{x_d} f)^{\top}$ for differentiable functions $f : \mathbb{R}^d \to \mathbb{R}$.

Standing Assumption 2

The distribution P admits a positive and differentiable PDF such that $x \mapsto (\nabla \log p)(x)$ is Lipschitz.

Definition 17 (Canonical Stein operator)

For a distribution P admitting a positive and differentiable density p on \mathbb{R}^d , we define the *canonical Stein operator*

$$(\mathcal{A}_{P}g)(x) = (\nabla \cdot g)(x) + g(x) \cdot (\nabla \log p)(x)$$

acting on differentiable vector field $g : \mathbb{R}^d \to \mathbb{R}^d$, where $\mathbf{x} \in \mathbb{R}^d$.

The canonical Stein operator was introduced (for Gaussian P) in Stein [1972]. Importantly, observe that

$$(\nabla \log p)(x) = \frac{(\nabla p)(x)}{p(x)} = \frac{\frac{1}{Z}(\nabla \tilde{p})(x)}{\frac{1}{Z}\tilde{p}(x)} = \frac{(\nabla \tilde{p})(x)}{\tilde{p}(x)} = (\nabla \log \tilde{p})(x),$$

which can be computed without knowledge of p or Z, provided \tilde{p} and $\nabla \tilde{p}$ can be evaluated.



Now we apply the Stein operator A_P to a standard kernel k, to obtain the following Stein kernel:

Lemma 18

Suppose the kernel $k: \mathbb{R}^d \times \mathbb{R}^d \to \mathbb{R}$ has $(x,y) \mapsto \partial^{(\alpha,\beta)} k(x,y)$ continuous and uniformly bounded for all $|\alpha|, |\beta| \leq 1$. Suppose $\int \|\nabla \log p(x)\| \mathrm{d} P(x) < \infty$ and that $\sup_{\|x\| \geq r} r^{d-1} p(x) \to 0$ as $r \to \infty$. Then

$$k_{P}(x,y) = \nabla_{x} \cdot \nabla_{y} k(x,y) + \nabla_{x} \log p(x) \cdot \nabla_{y} k(x,y) + \nabla_{y} \log p(y) \cdot \nabla_{x} k(x,y) + (\nabla_{x} \log p(x)) \cdot (\nabla_{y} \log p(y)) k(x,y)$$

is a kernel with $\int k_P(x,y) dP(y) = 0$ for all $x \in \mathbb{R}^d$.

Proof

First notice that

$$k_P(x,y) = \mathcal{A}_P^y \left[\begin{array}{c} \vdots \\ \nabla_{x_i} k(x,y) + k(x,y) \nabla_{x_i} \log p(x) \\ \vdots \end{array} \right] = \mathcal{A}_P^y g(y)$$

where, under our assumptions, (a) $y \mapsto g(y)$ is bounded, and (b) $y \mapsto \nabla_y \cdot g(y)$ is integrable with respect to P. Thus it suffices to show that $\int \mathcal{A}_P g dP = 0$ for all vector fields g for which (a) and (b) hold

Now we apply the Stein operator A_P to a standard kernel k, to obtain the following Stein kernel:

Lemma 18

Suppose the kernel $k: \mathbb{R}^d \times \mathbb{R}^d \to \mathbb{R}$ has $(x,y) \mapsto \partial^{(\alpha,\beta)} k(x,y)$ continuous and uniformly bounded for all $|\alpha|, |\beta| \leq 1$. Suppose $\int \|\nabla \log p(x)\| \mathrm{d} P(x) < \infty$ and that $\sup_{\|x\| \geq r} r^{d-1} p(x) \to 0$ as $r \to \infty$. Then

$$k_{P}(x,y) = \nabla_{x} \cdot \nabla_{y} k(x,y) + \nabla_{x} \log p(x) \cdot \nabla_{y} k(x,y) + \nabla_{y} \log p(y) \cdot \nabla_{x} k(x,y) + (\nabla_{x} \log p(x)) \cdot (\nabla_{y} \log p(y)) k(x,y)$$

is a kernel with $\int k_P(x,y) dP(y) = 0$ for all $x \in \mathbb{R}^d$.

Proof.

First notice that

$$k_P(\mathbf{x}, \mathbf{y}) = \mathcal{A}_P^{\mathbf{y}} \left[\begin{array}{c} \vdots \\ \nabla_{\mathbf{x}_i} k(\mathbf{x}, \mathbf{y}) + k(\mathbf{x}, \mathbf{y}) \nabla_{\mathbf{x}_i} \log p(\mathbf{x}) \\ \vdots \end{array} \right] = \mathcal{A}_P^{\mathbf{y}} g(\mathbf{y})$$

where, under our assumptions, (a) $y \mapsto g(y)$ is bounded, and (b) $y \mapsto \nabla_y \cdot g(y)$ is integrable with respect to P. Thus it suffices to show that $\int \mathcal{A}_P g dP = 0$ for all vector fields g for which (a) and (b) hold.

Now we apply the Stein operator A_P to a standard kernel k, to obtain the following Stein kernel:

Lemma 18

Suppose the kernel $k: \mathbb{R}^d \times \mathbb{R}^d \to \mathbb{R}$ has $(x,y) \mapsto \partial^{(\alpha,\beta)} k(x,y)$ continuous and uniformly bounded for all $|\alpha|, |\beta| \leq 1$. Suppose $\int \|\nabla \log p(x)\| dP(x) < \infty$ and that $\sup_{\|x\| \geq r} r^{d-1} p(x) \to 0$ as $r \to \infty$. Then

$$k_{P}(x,y) = \nabla_{x} \cdot \nabla_{y} k(x,y) + \nabla_{x} \log p(x) \cdot \nabla_{y} k(x,y) + \nabla_{y} \log p(y) \cdot \nabla_{x} k(x,y) + (\nabla_{x} \log p(x)) \cdot (\nabla_{y} \log p(y)) k(x,y)$$

is a kernel with $\int k_P(x,y) dP(y) = 0$ for all $x \in \mathbb{R}^d$.

Proof.

Let g be such a vector field, and let $B_r = \{x \in \mathbb{R}^d : ||x|| \le r\}$ and $S_r = \{x \in \mathbb{R}^d : ||x|| = r\}$. The main idea is to apply the divergence theorem (i.e. integrate by parts):

$$\int \mathcal{A}_{P}g dP = \int (\nabla \cdot g) + g \cdot (\nabla \log p) dP = \int (\nabla \cdot (pg))(x) dx$$
$$= \lim_{r \to \infty} \int_{\mathcal{B}_{r}} (\nabla \cdot (pg))(x) dx = \lim_{r \to \infty} \oint_{\mathcal{S}_{r}} p(x)(g(x) \cdot n(x)) dx$$

where n(x) is the outward unit normal to S_r at x. (The regularity assumptions ensure that the integrals $\int (\nabla \cdot g) dP$ and $\int g \cdot (\nabla \log p) dP$ exist.)

Now we apply the Stein operator A_P to a standard kernel k, to obtain the following Stein kernel:

Lemma 18

Suppose the kernel $k: \mathbb{R}^d \times \mathbb{R}^d \to \mathbb{R}$ has $(x,y) \mapsto \partial^{(\alpha,\beta)} k(x,y)$ continuous and uniformly bounded for all $|\alpha|, |\beta| \leq 1$. Suppose $\int \|\nabla \log p(x)\| \mathrm{d} P(x) < \infty$ and that $\sup_{\|x\| \geq r} r^{d-1} p(x) \to 0$ as $r \to \infty$. Then

$$k_{P}(x,y) = \nabla_{x} \cdot \nabla_{y} k(x,y) + \nabla_{x} \log p(x) \cdot \nabla_{y} k(x,y) + \nabla_{y} \log p(y) \cdot \nabla_{x} k(x,y) + (\nabla_{x} \log p(x)) \cdot (\nabla_{y} \log p(y)) k(x,y)$$

is a kernel with $\int k_P(x,y)dP(y) = 0$ for all $x \in \mathbb{R}^d$.

Proof.

Now

$$\oint_{S_r} p(x)(g(x) \cdot n(x)) dx \le ||g||_{\infty} \sup_{\|x\| \ge r} p(x) \oint_{S_r} dx$$

$$= ||g||_{\infty} \sup_{\|x\| \ge r} p(x) \frac{2\pi^{d/2}}{\Gamma(d/2)} r^{d-1}$$

$$\to 0 \text{ as } r \to \infty.$$

where we have used the formula for the surface area of the radius r sphere in \mathbb{R}^d .



For $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ and k_P defined in Lemma 18, we obtain the kernel Stein discrepancy (KSD):

$$D_{k_{\mathcal{P}}}(Q_n) = \sup_{\|f\|_{\mathcal{H}(k_{\mathcal{P}})} \leq 1} \left| \int f dQ_n \right| = \sqrt{\sum_{i=1}^n \sum_{j=1}^n w_i w_j k_{\mathcal{P}}(\mathbf{x}_i, \mathbf{x}_j)}$$

As with MMD, we can establish properties analogous to characteristicness and convergence control for KSD. Here we focus on convergence control:

Theorem 19 (Gorham and Mackey [2017])

Let P be distantly dissipative [see Gorham et al., 2019]. Consider the kernel

$$k(x,y) = (\sigma^2 + ||x - y||^2)^{-\beta}$$

for some fixed $\sigma > 0$ and a fixed exponent $\beta \in (0,1)$. Then $D_{k_P}(Q_n) \to 0$ implies $Q_n \Rightarrow P$.

For $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ and k_P defined in Lemma 18, we obtain the KSD:

$$D_{k_{\mathcal{P}}}(Q_n) = \sup_{\|f\|_{\mathcal{H}(k_{\mathcal{P}})} \leq 1} \left| \int f dQ_n \right| = \sqrt{\sum_{i=1}^n \sum_{j=1}^n w_i w_j k_{\mathcal{P}}(\mathbf{x}_i, \mathbf{x}_j)}$$

As with MMD, we can establish properties analogous to characteristicness and convergence control for KSD. Here we focus on convergence control:

Theorem 19 (Gorham and Mackey [2017])

Let P be distantly dissipative [see Gorham et al., 2019]. Consider the kernel

$$k(x, y) = (\sigma^{2} + ||x - y||^{2})^{-\beta}$$

for some fixed $\sigma > 0$ and a fixed exponent $\beta \in (0,1)$. Then $D_{k_P}(Q_n) \to 0$ implies $Q_n \Rightarrow P$.

For $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ and k_P defined in Lemma 18, we obtain the KSD:

$$D_{k_{\mathcal{P}}}(Q_n) = \sup_{\|f\|_{\mathcal{H}(k_{\mathcal{P}})} \leq 1} \left| \int f dQ_n \right| = \sqrt{\sum_{i=1}^n \sum_{j=1}^n w_i w_j k_{\mathcal{P}}(\mathbf{x}_i, \mathbf{x}_j)}$$

As with MMD, we can establish properties analogous to characteristicness and convergence control for KSD. Here we focus on convergence control:

Theorem 19 (Gorham and Mackey [2017])

Let P be distantly dissipative [see Gorham et al., 2019]. Consider the kernel

$$k(\mathbf{x}, \mathbf{y}) = (\sigma^2 + \|\mathbf{x} - \mathbf{y}\|^2)^{-\beta}$$

for some fixed $\sigma > 0$ and a fixed exponent $\beta \in (0,1)$. Then $D_{kp}(Q_n) \to 0$ implies $Q_n \Rightarrow P$.

For $Q_n = \sum_{i=1}^n w_i \delta(\mathbf{x}_i)$ and k_P defined in Lemma 18, we obtain the KSD:

$$D_{k_{\mathcal{P}}}(Q_n) = \sup_{\|f\|_{\mathcal{H}(k_{\mathcal{P}})} \le 1} \left| \int f dQ_n \right| = \sqrt{\sum_{i=1}^n \sum_{j=1}^n w_i w_j k_{\mathcal{P}}(\mathbf{x}_i, \mathbf{x}_j)}$$

As with MMD, we can establish properties analogous to characteristicness and convergence control for KSD. Here we focus on convergence control:

Theorem 19 (Gorham and Mackey [2017])

Let P be distantly dissipative [see Gorham et al., 2019]. Consider the kernel

$$k(x, y) = (\sigma^2 + ||x - y||^2)^{-\beta}$$

for some fixed $\sigma > 0$ and a fixed exponent $\beta \in (0,1)$. Then $D_{kp}(Q_n) \to 0$ implies $Q_n \Rightarrow P$.



Optimal Quantisation

The algorithms previously described for MMD can immediately be applied to KSD!

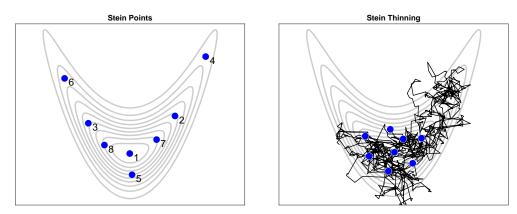


Figure: Stein points (left) and Stein thinning (right): Stein points are generated by sequential (greedy) minimisation of KSD to select n=8 states $\{x_i\}_{i=1}^n$ in an approximation Q_n to P. The numbers indicate the order in which the states x_i were selected. Stein thinning restricts the continuous inner-loop optimisation problem in Stein points to a discrete search over a MCMC sample path (black).

Optimal weights can again be computed, with a slight twist:

Lemma 20

Let $x_1, \ldots, x_n \in \mathcal{X}$ be distinct. The optimal weights

$$\underset{\substack{w \in \mathbb{R}^n \\ 1^\top w = 1}}{\arg \min} D_{kp} \left(\sum_{i=1}^n w_i \delta(x_i) \right)$$

are $w = (K_P^{-1}1)/(1^T K_P^{-1}1)$, where $[K_P]_{ij} = k_P(x_i, x_j)$. [Note the constraint $w_1 + \cdots + w_n = 1$.]

Proof

From the explicit form of MMD we have

$$D_{k_P}\left(\sum_{i=1}^n w_i \delta(x_i)\right)^2 = w^\top K_P w,$$

so the optimisation problem i

$$\operatorname{arg\,min} \mathbf{w}^{\top} K_P \mathbf{w}$$
 s.t. $\mathbf{1}^{\top} \mathbf{w} = 1$.

This can be solved using the method of Lagrange multipliers to obtain the stated result.

Optimal weights can again be computed, with a slight twist:

Lemma 20

Let $x_1, \ldots, x_n \in \mathcal{X}$ be distinct. The optimal weights

$$\underset{\substack{\mathbf{w} \in \mathbb{R}^n \\ \mathbf{1}^{\top} \mathbf{w} = 1}}{\arg \min} D_{k_P} \left(\sum_{i=1}^n w_i \delta(\mathbf{x}_i) \right)$$

are
$$\mathbf{w} = (K_P^{-1}\mathbf{1})/(\mathbf{1}^{\top}K_P^{-1}\mathbf{1})$$
, where $[K_P]_{ij} = k_P(\mathbf{x}_i, \mathbf{x}_j)$. [Note the constraint $w_1 + \dots + w_n = 1$.]

Proof

From the explicit form of MMD we have

$$D_{k_P}\left(\sum_{i=1}^n w_i \delta(x_i)\right)^2 = \mathbf{w}^\top \mathbf{K}_P \mathbf{w},$$

so the optimisation problem is

$$\operatorname{arg\,min\,} w^{\top} K_P w$$
 s.t. $\mathbf{1}^{\top} w = 1$.

This can be solved using the method of Lagrange multipliers to obtain the stated result.



Optimal weights can again be computed, with a slight twist:

Lemma 20

Let $x_1, \ldots, x_n \in \mathcal{X}$ be distinct. The optimal weights

$$\underset{\substack{\mathbf{w} \in \mathbb{R}^n \\ \mathbf{1}^{\top} \mathbf{w} = 1}}{\arg \min} D_{k_P} \left(\sum_{i=1}^n w_i \delta(\mathbf{x}_i) \right)$$

are
$$\mathbf{w} = (K_P^{-1}\mathbf{1})/(\mathbf{1}^{\top}K_P^{-1}\mathbf{1})$$
, where $[K_P]_{ij} = k_P(\mathbf{x}_i, \mathbf{x}_j)$. [Note the constraint $w_1 + \cdots + w_n = 1$.]

Proof

From the explicit form of MMD we have

$$D_{k_P}\left(\sum_{i=1}^n w_i \delta(x_i)\right)^2 = \mathbf{w}^\top \mathbf{K}_P \mathbf{w},$$

so the optimisation problem is

$$\operatorname{arg\,min} w^{\top} K_P w$$
 s.t. $\mathbf{1}^{\top} w = 1$.

This can be solved using the method of Lagrange multipliers to obtain the stated result

Optimal weights can again be computed, with a slight twist:

Lemma 20

Let $x_1, \ldots, x_n \in \mathcal{X}$ be distinct. The optimal weights

$$\underset{\substack{\mathbf{w} \in \mathbb{R}^n \\ \mathbf{1}^{\top} \mathbf{w} = 1}}{\operatorname{arg \, min}} D_{k_{P}} \left(\sum_{i=1}^{n} w_{i} \delta(\mathbf{x}_{i}) \right)$$

are $\mathbf{w} = (K_P^{-1}\mathbf{1})/(\mathbf{1}^{\top}K_P^{-1}\mathbf{1})$, where $[K_P]_{ij} = k_P(\mathbf{x}_i, \mathbf{x}_j)$. [Note the constraint $w_1 + \cdots + w_n = 1$.]

Proof.

From the explicit form of MMD we have

$$D_{k_P}\left(\sum_{i=1}^n w_i \delta(\mathbf{x}_i)\right)^2 = \mathbf{w}^\top \mathbf{K}_P \mathbf{w},$$

so the optimisation problem is

$$\operatorname{arg\,min} \mathbf{w}^{\top} \mathbf{K}_{P} \mathbf{w}$$
 s.t. $\mathbf{1}^{\top} \mathbf{w} = 1$.

This can be solved using the method of Lagrange multipliers to obtain the stated result.



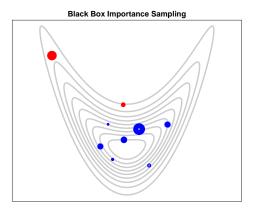


Figure: Black box importance sampling [Liu and Lee, 2017, Hodgkinson et al., 2020]: In black box importance sampling the weights w_1, \ldots, w_n are obtained by minimising KSD in the manner of Lemma 20. Blue indicates states x_i with positive weights $w_i > 0$, while red indicates negative weights $w_i < 0$. The size of the circles is proportional to $|w_i|$.

Complexity = $O(n^3)$ and symmetry exploits fail, but $P_{ST} \to P_{BBIS}$ as $m \to \infty$ for n fixed



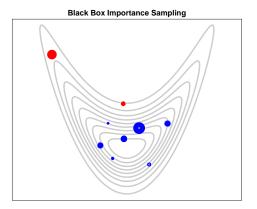
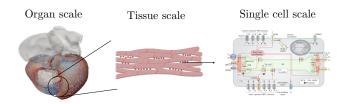


Figure: Black box importance sampling [Liu and Lee, 2017, Hodgkinson et al., 2020]: In black box importance sampling the weights w_1, \ldots, w_n are obtained by minimising KSD in the manner of Lemma 20. Blue indicates states x_i with positive weights $w_i > 0$, while red indicates negative weights $w_i < 0$. The size of the circles is proportional to $|w_i|$.

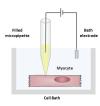
Complexity = $O(n^3)$ and symmetry exploits fail, but $P_{ST} \to P_{BBIS}$ as $m \to \infty$ for n fixed.

Case Study: Cardiac Digital Twins

Cardiac Models

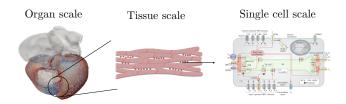


- Cardiac function determined by integrated action of myocytes.
- Calcium is intracellular end-point signal driving contraction.
- **Question**: is there cell-to-cell variability?
 - "Extreme" cells might be related to pathologies (arrythmias ventricular fibrillation)?
 - ▶ Role for ageing in the extent of heterogeneity?
 - **•** . . .
- 25 myocytes, from different rat hearts, fed in a dish
- Patch clamp experiment: micropipette controls cell membrane potential, and adds a drug to
 - block transients of ions that are not calcium:
 - stimulate activity of calcium handling proteins (facilitate system identification).

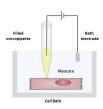




Cardiac Models

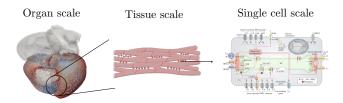


- Cardiac function determined by integrated action of myocytes.
- Calcium is intracellular end-point signal driving contraction.
- ▶ Question: is there cell-to-cell variability?
 - Extreme" cells might be related to pathologies (arrythmias ventricular fibrillation)?
 - Role for ageing in the extent of heterogeneity?
 - ▶ ...
- 25 myocytes, from different rat hearts, fed in a dish.
- Patch clamp experiment: micropipette controls cell membrane potential, and adds a drug to
 - block transients of ions that are not calcium;
 - stimulate activity of calcium handling proteins (facilitate system identification).

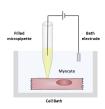




Cardiac Models



- Cardiac function determined by integrated action of myocytes.
- Calcium is intracellular end-point signal driving contraction.
- ▶ Question: is there cell-to-cell variability?
 - "Extreme" cells might be related to pathologies (arrythmias ventricular fibrillation)?
 - Role for ageing in the extent of heterogeneity?
 - **.**..
- 25 myocytes, from different rat hearts, fed in a dish
- Patch clamp experiment: micropipette controls cell membrane potential, and adds a drug to
 - block transients of ions that are not calcium;
 - stimulate activity of calcium handling proteins (facilitate system identification).





Statistical Model

- ▶ The experiment is modelled using mass action kinetics, in the form of an ordinary differential equation (ODE) with unknown parameters denoted $\theta \in \Theta$.
- ► The generic form of the ODE is

$$u(0) = u_0, \qquad \frac{\mathrm{d}u_1}{\mathrm{d}t} = f_1(t, u; \theta)$$

$$\vdots$$

$$\frac{\mathrm{d}u_p}{\mathrm{d}t} = f_p(t, u; \theta)$$

and the solution at time t is denoted $u(t;\theta)$. Here p=6 and $d:=\dim(\Theta)=38$

The data are related to the ODE via a measurement error mode

$$\pi(y|\theta) := \prod_{i=1}^n \frac{1}{\sqrt{2\pi}\sigma} \exp\left(-\frac{(y_i - u_1(t_i;\theta))^2}{2\sigma^2}\right)$$

where σ is the resolution of the measurement equipment. (Less trivial measurement models may be needed, including temporal components in the measurement error.)

(We assume u_0 and σ are known in this talk.)

Statistical Model

- ▶ The experiment is modelled using mass action kinetics, in the form of an ordinary differential equation (ODE) with unknown parameters denoted $\theta \in \Theta$.
- ▶ The generic form of the ODE is

$$u(0) = u_0,$$
 $\frac{\mathrm{d}u_1}{\mathrm{d}t} = f_1(t, u; \theta),$ \vdots $\frac{\mathrm{d}u_p}{\mathrm{d}t} = f_p(t, u; \theta),$

and the solution at time t is denoted $u(t; \theta)$. Here p = 6 and $d := \dim(\Theta) = 38$.

The data are related to the ODE via a measurement error mode

$$\pi(y|\theta) := \prod_{i=1}^n rac{1}{\sqrt{2\pi}\sigma} \exp\left(-rac{(y_i - u_1(t_i; heta))^2}{2\sigma^2}
ight),$$

where σ is the resolution of the measurement equipment. (Less trivial measurement models may be needed, including temporal components in the measurement error.)

(We assume u_0 and σ are known in this talk.)



Statistical Model

- ▶ The experiment is modelled using mass action kinetics, in the form of an ordinary differential equation (ODE) with unknown parameters denoted $\theta \in \Theta$.
- ► The generic form of the ODE is

$$u(0) = u_0,$$
 $\frac{\mathrm{d}u_1}{\mathrm{d}t} = f_1(t, u; \theta),$ \vdots $\frac{\mathrm{d}u_{\rho}}{\mathrm{d}t} = f_{\rho}(t, u; \theta),$

and the solution at time t is denoted $u(t;\theta)$. Here p=6 and $d:=\dim(\Theta)=38$.

▶ The data are related to the ODE via a measurement error model

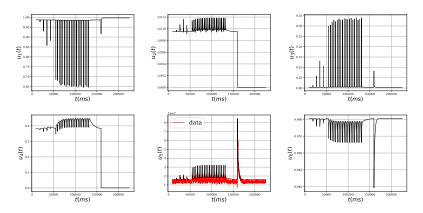
$$\pi(y|\theta) := \prod_{i=1}^n \frac{1}{\sqrt{2\pi}\sigma} \exp\left(-\frac{(y_i - u_1(t_i;\theta))^2}{2\sigma^2}\right),$$

where σ is the resolution of the measurement equipment. (Less trivial measurement models may be needed, including temporal components in the measurement error.)

(We assume u_0 and σ are known in this talk.)

Inverse Problem

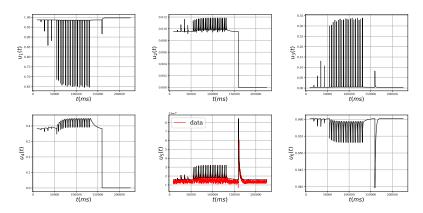
- Is cell-to-cell heterogeneity observed at the level of the parameters θ ? (The figure corresponds to slightly perturbed literature values for θ).
- ightharpoonup To answer this we need to solve an inverse problem for θ in the statistical model.



► Challenge: ODE solver failure.

Inverse Problem

- Is cell-to-cell heterogeneity observed at the level of the parameters θ ? (The figure corresponds to slightly perturbed literature values for θ).
- ightharpoonup To answer this we need to solve an inverse problem for θ in the statistical model.



Challenge: ODE solver failure.

MCMC samples are biased by ODE solver failure, but Stein Thinning does not require MCMC to be P-invariant - as long as the relevant part of the parameter space is explored:

Theorem 21 ([Riabiz et al., 2022])

Let $(\theta_i)_{i \in \mathbb{N}}$ be a Q-invariant, time-homogeneous, reversible Markov chain, such that P is absolutely continuous with respect to Q and

- $(\theta_i)_{i\in\mathbb{N}}$ is V-uniformly ergodic with $V(\theta) \geq \frac{dP}{dQ}(\theta)\sqrt{k_P(\theta,\theta)}$
- $ightharpoonup \sup_{i\in\mathbb{N}}\mathbb{E}[\frac{\mathrm{d}P}{\mathrm{d}Q}(\theta_i)\sqrt{k_P(\theta_i,\theta_i)}V(\theta_i)]<\infty$

Then the output of Stein Thinning satisfies

$$\frac{1}{m}\sum_{i\in S}\delta(\theta_i)\Rightarrow P$$

MCMC samples are biased by ODE solver failure, but Stein Thinning does not require MCMC to be P-invariant - as long as the relevant part of the parameter space is explored:

Theorem 21 ([Riabiz et al., 2022])

Let $(\theta_i)_{i\in\mathbb{N}}$ be a Q-invariant, time-homogeneous, reversible Markov chain, such that P is absolutely continuous with respect to Q and

- $(\theta_i)_{i\in\mathbb{N}}$ is V-uniformly ergodic with $V(\theta) \geq \frac{dP}{dQ}(\theta)\sqrt{k_P(\theta,\theta)}$
- $ightharpoonup \sup_{i\in\mathbb{N}}\mathbb{E}[\frac{\mathrm{d}P}{\mathrm{d}Q}(\theta_i)\sqrt{k_P(\theta_i,\theta_i)}V(\theta_i)]<\infty$
- $\qquad \exists \gamma > 0 \text{ s.t. } b := \sup_{i \in \mathbb{N}} \mathbb{E}[e^{\gamma \max(1, \frac{\mathrm{d}P}{\mathrm{d}Q}(\theta_i)^2)k_P(\theta_i, \theta_i)}] < \infty.$

Then the output of Stein Thinning satisfies

$$\frac{1}{m}\sum_{i\in S}\delta(\theta_i)\Rightarrow P$$

MCMC samples are biased by ODE solver failure, but Stein Thinning does not require MCMC to be P-invariant - as long as the relevant part of the parameter space is explored:

Theorem 21 ([Riabiz et al., 2022])

Let $(\theta_i)_{i\in\mathbb{N}}$ be a Q-invariant, time-homogeneous, reversible Markov chain, such that P is absolutely continuous with respect to Q and

- $(\theta_i)_{i\in\mathbb{N}}$ is V-uniformly ergodic with $V(\theta) \geq \frac{dP}{dQ}(\theta)\sqrt{k_P(\theta,\theta)}$
- $ightharpoonup \sup_{i\in\mathbb{N}}\mathbb{E}[\frac{\mathrm{d}P}{\mathrm{d}Q}(\theta_i)\sqrt{k_P(\theta_i,\theta_i)}V(\theta_i)]<\infty$
- $\qquad \exists \gamma > 0 \text{ s.t. } b := \sup_{i \in \mathbb{N}} \mathbb{E}[e^{\gamma \max(1, \frac{\mathrm{d} P}{\mathrm{d} Q}(\theta_i)^2)k_P(\theta_i, \theta_i)}] < \infty.$

Then the output of Stein Thinning satisfies

$$\frac{1}{m}\sum_{i\in S}\delta(\theta_i)\Rightarrow P$$

MCMC samples are biased by ODE solver failure, but Stein Thinning does not require MCMC to be P-invariant - as long as the relevant part of the parameter space is explored:

Theorem 21 ([Riabiz et al., 2022])

Let $(\theta_i)_{i\in\mathbb{N}}$ be a Q-invariant, time-homogeneous, reversible Markov chain, such that P is absolutely continuous with respect to Q and

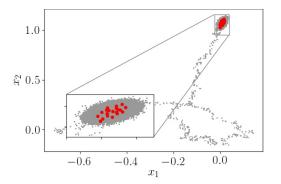
- $(\theta_i)_{i\in\mathbb{N}}$ is V-uniformly ergodic with $V(\theta) \geq \frac{dP}{dQ}(\theta)\sqrt{k_P(\theta,\theta)}$
- $ightharpoonup \sup_{i\in\mathbb{N}}\mathbb{E}[\frac{\mathrm{d}P}{\mathrm{d}Q}(\theta_i)\sqrt{k_P(\theta_i,\theta_i)}V(\theta_i)]<\infty$
- $\qquad \exists \gamma > 0 \text{ s.t. } b := \sup_{i \in \mathbb{N}} \mathbb{E}[e^{\gamma \max(1, \frac{dP}{dQ}(\theta_i)^2)k_P(\theta_i, \theta_i)}] < \infty.$

Then the output of Stein Thinning satisfies

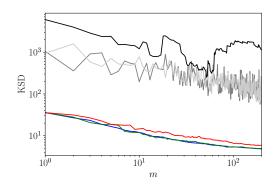
$$\frac{1}{m}\sum_{i\in S}\delta(\theta_i)\Rightarrow P$$

Illustrative Application to Differential Equation Constrained Inverse Problems

Goodwin oscillator; d=4 parameters to be estimated. (Red dots are Stein Thinning, while gray dots are MCMC.)



Cardiac model; d=38 parameters to be esitmated. (Blue, red, and green are Stein Thinning, while black are MCMC.)



Future Directions, Open Questions and Challenges

Scalable Stein Thinning

Greedy selection may be sub-optimal. Also, the cost of selecting m points from n using Stein Thinning is high, at $O(m^2n)$.

- A **non-myopic** algorithm selects *s* points simultaneously.
- A mini-batch algorithm searches over a subset of $b \ll n$ candidates at each step.

Full details in Teymur et al. [2021]

Scalable Stein Thinning

Greedy selection may be sub-optimal. Also, the cost of selecting m points from n using Stein Thinning is high, at $O(m^2n)$.

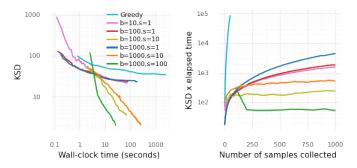
- A **non-myopic** algorithm selects *s* points simultaneously.
- ▶ A mini-batch algorithm searches over a subset of $b \ll n$ candidates at each step.

Full details in Teymur et al. [2021]

Scalable Stein Thinning

Greedy selection may be sub-optimal. Also, the cost of selecting m points from n using Stein Thinning is high, at $O(m^2n)$.

- A non-myopic algorithm selects s points simultaneously.
- ▶ A **mini-batch** algorithm searches over a subset of $b \ll n$ candidates at each step.



Full details in Teymur et al. [2021].

Gradients can be difficult to stably compute, e.g. for ODEs and PDEs.

Definition 22 (Gradent-free Stein operator)

For distributions P and Q admitting sufficiently regular densities p and q on \mathbb{R}^d , we define the gradient-free Stein operator

$$(\mathcal{A}_{P,Q} g)(x) = \frac{q(x)}{p(x)} [(\nabla \cdot g)(x) + g(x) \cdot (\nabla \log q)(x)]$$

- ► The associated Stein discrepancy can be computed up to proportionality using $\tilde{p}(x) = Z \times p(x)$.
- ► The differentiability and distant dissipativity requirements now apply to q (a degree of freedom) instead of p (the posterior target).
- Sufficient conditions have been established on q for convergence detection and control.

Gradients can be difficult to stably compute, e.g. for ODEs and PDEs.

Definition 22 (Gradent-free Stein operator)

For distributions P and Q admitting sufficiently regular densities p and q on \mathbb{R}^d , we define the gradient-free Stein operator

$$(\mathcal{A}_{P,Q} g)(x) = \frac{q(x)}{p(x)} [(\nabla \cdot g)(x) + g(x) \cdot (\nabla \log q)(x)]$$

- The associated Stein discrepancy can be computed up to proportionality using $\tilde{p}(x) = Z \times p(x)$.
- ► The differentiability and distant dissipativity requirements now apply to *q* (a degree of freedom) instead of *p* (the posterior target).
- Sufficient conditions have been established on q for convergence detection and control.

Gradients can be difficult to stably compute, e.g. for ODEs and PDEs.

Definition 22 (Gradent-free Stein operator)

For distributions P and Q admitting sufficiently regular densities p and q on \mathbb{R}^d , we define the gradient-free Stein operator

$$(\mathcal{A}_{P,Q} g)(x) = \frac{q(x)}{p(x)} [(\nabla \cdot g)(x) + g(x) \cdot (\nabla \log q)(x)]$$

- ► The associated Stein discrepancy can be computed up to proportionality using $\tilde{p}(x) = Z \times p(x)$.
- ► The differentiability and distant dissipativity requirements now apply to *q* (a degree of freedom) instead of *p* (the posterior target).
- Sufficient conditions have been established on q for convergence detection and control.

Gradients can be difficult to stably compute, e.g. for ODEs and PDEs.

Definition 22 (Gradent-free Stein operator)

For distributions P and Q admitting sufficiently regular densities p and q on \mathbb{R}^d , we define the gradient-free Stein operator

$$(\mathcal{A}_{P,Q} g)(x) = \frac{q(x)}{p(x)} [(\nabla \cdot g)(x) + g(x) \cdot (\nabla \log q)(x)]$$

- The associated Stein discrepancy can be computed up to proportionality using $\tilde{p}(x) = Z \times p(x)$.
- ► The differentiability and distant dissipativity requirements now apply to *q* (a degree of freedom) instead of *p* (the posterior target).
- Sufficient conditions have been established on a for convergence detection and control.

Gradients can be difficult to stably compute, e.g. for ODEs and PDEs.

Definition 22 (Gradent-free Stein operator)

For distributions P and Q admitting sufficiently regular densities p and q on \mathbb{R}^d , we define the gradient-free Stein operator

$$(\mathcal{A}_{P,Q} g)(x) = \frac{q(x)}{p(x)} [(\nabla \cdot g)(x) + g(x) \cdot (\nabla \log q)(x)]$$

- The associated Stein discrepancy can be computed up to proportionality using $\tilde{p}(x) = Z \times p(x)$.
- ► The differentiability and distant dissipativity requirements now apply to q (a degree of freedom) instead of p (the posterior target).
- Sufficient conditions have been established on q for convergence detection and control.

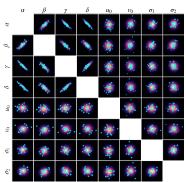
Gradients can be difficult to stably compute, e.g. for ODEs and PDEs.

Definition 22 (Gradent-free Stein operator)

For distributions P and Q admitting sufficiently regular densities p and q on \mathbb{R}^d , we define the gradient-free Stein operator

$$(A_{P,Q} g)(x) = \frac{q(x)}{p(x)} [(\nabla \cdot g)(x) + g(x) \cdot (\nabla \log q)(x)]$$

- The associated Stein discrepancy can be computed up to proportionality using $\tilde{p}(x) = Z \times p(x)$.
- ► The differentiability and distant dissipativity requirements now apply to q (a degree of freedom) instead of p (the posterior target).
- Sufficient conditions have been established on q for convergence detection and control.



Software Development

To try out these algorithms (in Python, Matlab and R), visit Stein-Thinning.org

Stein Thinning About Stein Thinning is a tool for post-processing the output of a sampling procedure, such as Markov chain Monte Carlo (MCMC). It aims to minimise a Stein discrepancy, selecting a subsequence of samples that best represent the distributional target. J.ein Thinning **Bias Correction for MCMC** Optimally thinning of output from a sampling procedure, such as MCMC. Here the red samples are automatically chosen by Stein Thinning to provide a more accurate approximation to the distributional target, compared with the original MCMC output. [Read more] Watch on YouTube View the Project on GitHub The user provides two arrays: one containing the samples and another containing the corresponding gradients of the log-target. Stein Thinning returns a vector of indices, indicating which samples were selected. In favourable circumstances, Stein Thinning is able to: · automatically identify and remove the burn-in period from MCMC, · perform bias-removal for biased sampling procedures, · provide improved approximations of the distributional target, · offer a compressed representation of sample-based output.

Broader Context: Optimisation over $\mathcal{P}(\Theta)$

Going beyond optimisation in Θ , we can consider optimisation in $\mathcal{P}(\Theta)$:

$$D_{\mathcal{H}(k_P),P}(Q) := \sup_{\|f\|_{\mathcal{H}(k_P)} \leq 1} |\mathbb{E}_{\vartheta \sim Q}[f(\vartheta)]|$$

▶ Variational Inference: Ranganath et al. [2016], Hu et al. [2018], Fisher et al. [2021], ...

$$\min_{Q\in\mathcal{Q}} D_{\mathcal{H}(k_P),P}(Q), \qquad \text{(e.g.)} \ \ \mathcal{Q} = \{ T_\# Q_0 : T \text{ a neural network } \}$$

Avoids the requirement in VI that T be a diffeomorphism (i.e. no need for normalising flows!)

► **Gradient Flow:** Korba et al. [2021]

$$\frac{\partial Q_t}{\partial t} + \operatorname{div}(Q_t v_{Q_t}) = 0, \qquad v_{Q_t} = -\nabla_{W_2} \mathcal{F}(Q_t), \qquad \mathcal{F}(Q) = \frac{1}{2} D_{\mathcal{H}(k_P), P}(Q)^2$$

stnot the same as SVGD [see Liu, 2017].

Broader Context: Optimisation over $\mathcal{P}(\Theta)$

Going beyond optimisation in Θ , we can consider optimisation in $\mathcal{P}(\Theta)$:

$$D_{\mathcal{H}(k_P),P}(Q) := \sup_{\|f\|_{\mathcal{H}(k_P)} \le 1} |\mathbb{E}_{\vartheta \sim Q}[f(\vartheta)]|$$

▶ Variational Inference: Ranganath et al. [2016], Hu et al. [2018], Fisher et al. [2021], ...

$$\min_{Q\in\mathcal{Q}} D_{\mathcal{H}(k_P),P}(Q), \qquad \text{(e.g.)} \ \ \mathcal{Q} = \{T_\# Q_0 : T \text{ a neural network}\}$$

Avoids the requirement in VI that T be a diffeomorphism (i.e. no need for normalising flows!).

► **Gradient Flow:** Korba et al. [2021]

$$\frac{\partial Q_t}{\partial t} + \operatorname{div}(Q_t v_{Q_t}) = 0, \qquad v_{Q_t} = -\nabla_{W_2} \mathcal{F}(Q_t), \qquad \mathcal{F}(Q) = \frac{1}{2} D_{\mathcal{H}(k_P), P}(Q)^2$$

*not the same as SVGD [see Liu, 2017].

Broader Context: Optimisation over $\mathcal{P}(\Theta)$

Going beyond optimisation in Θ , we can consider optimisation in $\mathcal{P}(\Theta)$:

$$D_{\mathcal{H}(k_P),P}(Q) := \sup_{\|f\|_{\mathcal{H}(k_P)} \leq 1} |\mathbb{E}_{\vartheta \sim Q}[f(\vartheta)]|$$

▶ Variational Inference: Ranganath et al. [2016], Hu et al. [2018], Fisher et al. [2021], ...

$$\min_{Q\in\mathcal{Q}}D_{\mathcal{H}(k_P),P}(Q), \qquad \text{(e.g.)} \ \ \mathcal{Q}=\{\mathit{T}_\#\mathit{Q}_0: \mathit{T} \text{ a neural network}\}$$

Avoids the requirement in VI that T be a diffeomorphism (i.e. no need for normalising flows!).

► Gradient Flow: Korba et al. [2021]

$$\frac{\partial Q_t}{\partial t} + \text{div}(Q_t v_{Q_t}) = 0, \qquad v_{Q_t} = -\nabla_{W_2} \mathcal{F}(Q_t), \qquad \mathcal{F}(Q) = \frac{1}{2} D_{\mathcal{H}(k_P),P}(Q)^2$$

*not the same as SVGD [see Liu, 2017].



For any Stein characterisation (A, \mathcal{F}) we can consider an associated Stein discrepancy [Gorham and Mackey, 2015]:

$$D_{\mathcal{H}(k_P),P}(Q) := \sup_{f \in \mathcal{F}} |\mathbb{E}_{\vartheta \sim Q}[f(\vartheta)]|$$

- Beyond Euclidean State Spaces: Riemannian manifolds [Barp et al., 2022, Le et al., 2020]
 discrete spaces [Xu and Reinert, 2021], ...
- Beyond Kernel Stein Sets: bounded Lipschitz [Gorham and Mackey, 2015], neural network [Grathwohl et al., 2020], ...
- ▶ Beyond the Canonical Stein Operator: diffusion Stein operators [Gorham et al., 2019], ..
- Scalable Stein Discrepancies: random features [Huggins and Mackey, 2018], data sub-sampling [Gorham et al., 2020], slicing [Gong et al., 2020], ...

For any Stein characterisation (A, \mathcal{F}) we can consider an associated Stein discrepancy [Gorham and Mackey, 2015]:

$$D_{\mathcal{H}(k_P),P}(Q) := \sup_{f \in \mathcal{F}} |\mathbb{E}_{\vartheta \sim Q}[f(\vartheta)]|$$

- Beyond Euclidean State Spaces: Riemannian manifolds [Barp et al., 2022, Le et al., 2020], discrete spaces [Xu and Reinert, 2021], ...
- Beyond Kernel Stein Sets: bounded Lipschitz [Gorham and Mackey, 2015], neural network [Grathwohl et al., 2020], ...
- ▶ Beyond the Canonical Stein Operator: diffusion Stein operators [Gorham et al., 2019], ..
- Scalable Stein Discrepancies: random features [Huggins and Mackey, 2018], data sub-sampling [Gorham et al., 2020], slicing [Gong et al., 2020], ...

For any Stein characterisation (A, \mathcal{F}) we can consider an associated Stein discrepancy [Gorham and Mackey, 2015]:

$$D_{\mathcal{H}(k_P),P}(Q) := \sup_{f \in \mathcal{F}} |\mathbb{E}_{\vartheta \sim Q}[f(\vartheta)]|$$

- Beyond Euclidean State Spaces: Riemannian manifolds [Barp et al., 2022, Le et al., 2020], discrete spaces [Xu and Reinert, 2021], ...
- Beyond Kernel Stein Sets: bounded Lipschitz [Gorham and Mackey, 2015], neural network [Grathwohl et al., 2020], ...
- ▶ Beyond the Canonical Stein Operator: diffusion Stein operators [Gorham et al., 2019], ...
- Scalable Stein Discrepancies: random features [Huggins and Mackey, 2018], data sub-sampling [Gorham et al., 2020], slicing [Gong et al., 2020], ...

For any Stein characterisation (A, \mathcal{F}) we can consider an associated Stein discrepancy [Gorham and Mackey, 2015]:

$$D_{\mathcal{H}(k_P),P}(Q) := \sup_{f \in \mathcal{F}} |\mathbb{E}_{\vartheta \sim Q}[f(\vartheta)]|$$

- Beyond Euclidean State Spaces: Riemannian manifolds [Barp et al., 2022, Le et al., 2020], discrete spaces [Xu and Reinert, 2021], ...
- Beyond Kernel Stein Sets: bounded Lipschitz [Gorham and Mackey, 2015], neural network [Grathwohl et al., 2020], ...
- Beyond the Canonical Stein Operator: diffusion Stein operators [Gorham et al., 2019], ...
- Scalable Stein Discrepancies: random features [Huggins and Mackey, 2018], data sub-sampling [Gorham et al., 2020], slicing [Gong et al., 2020], ...

For any Stein characterisation (A, \mathcal{F}) we can consider an associated Stein discrepancy [Gorham and Mackey, 2015]:

$$D_{\mathcal{H}(k_P),P}(Q) := \sup_{f \in \mathcal{F}} |\mathbb{E}_{\vartheta \sim Q}[f(\vartheta)]|$$

- Beyond Euclidean State Spaces: Riemannian manifolds [Barp et al., 2022, Le et al., 2020], discrete spaces [Xu and Reinert, 2021], ...
- Beyond Kernel Stein Sets: bounded Lipschitz [Gorham and Mackey, 2015], neural network [Grathwohl et al., 2020], ...
- Beyond the Canonical Stein Operator: diffusion Stein operators [Gorham et al., 2019], ...
- Scalable Stein Discrepancies: random features [Huggins and Mackey, 2018], data sub-sampling [Gorham et al., 2020], slicing [Gong et al., 2020], ...

For any Stein characterisation (A, \mathcal{F}) we can consider an associated Stein discrepancy [Gorham and Mackey, 2015]:

$$D_{\mathcal{H}(k_P),P}(Q) := \sup_{f \in \mathcal{F}} |\mathbb{E}_{\vartheta \sim Q}[f(\vartheta)]|$$

- Beyond Euclidean State Spaces: Riemannian manifolds [Barp et al., 2022, Le et al., 2020], discrete spaces [Xu and Reinert, 2021], ...
- Beyond Kernel Stein Sets: bounded Lipschitz [Gorham and Mackey, 2015], neural network [Grathwohl et al., 2020], ...
- Beyond the Canonical Stein Operator: diffusion Stein operators [Gorham et al., 2019], ...
- Scalable Stein Discrepancies: random features [Huggins and Mackey, 2018], data sub-sampling [Gorham et al., 2020], slicing [Gong et al., 2020], ...

Broader Context: Alternatives to Direct Minimisation of Stein Discrepancy

- Stein Variational Gradient Descent: Liu and Wang [2016], Liu [2017], Liu and Zhu [2018], Detommaso et al. [2018], ...
- ▶ MCMC with Stein Control Variates: Assaraf and Caffarel [1999], Mira et al. [2013], CJO et al. [2017], Belomestny et al. [2017], South et al. [2022], ...

Given a Qol
$$f$$
, seek (u,c) such that $c+\frac{\nabla\cdot(p\nabla u)}{p}=f$. Then $c=\mathbb{E}_{\vartheta\sim P}[f(\vartheta)]$

In practice, an approximate solution u gives rise to a control variate $v = \nabla \cdot (p\nabla u)/p$ for use in MCMC.

For more, please see the recent review paper of Anastasiou et al. [2022].

Broader Context: Alternatives to Direct Minimisation of Stein Discrepancy

- Stein Variational Gradient Descent: Liu and Wang [2016], Liu [2017], Liu and Zhu [2018], Detommaso et al. [2018], ...
- ▶ MCMC with Stein Control Variates: Assaraf and Caffarel [1999], Mira et al. [2013], CJO et al. [2017], Belomestny et al. [2017], South et al. [2022], ...

Given a Qol
$$f$$
, seek (u,c) such that $c+\frac{\nabla\cdot(p\nabla u)}{p}=f$. Then $c=\mathbb{E}_{\vartheta\sim P}[f(\vartheta)]$.

In practice, an approximate solution u gives rise to a control variate $v = \nabla \cdot (p\nabla u)/p$ for use in MCMC.

For more, please see the recent review paper of Anastasiou et al. [2022].

Broader Context: Alternatives to Direct Minimisation of Stein Discrepancy

- Stein Variational Gradient Descent: Liu and Wang [2016], Liu [2017], Liu and Zhu [2018], Detommaso et al. [2018], ...
- ▶ MCMC with Stein Control Variates: Assaraf and Caffarel [1999], Mira et al. [2013], CJO et al. [2017], Belomestny et al. [2017], South et al. [2022], ...

Given a Qol
$$f$$
, seek (u,c) such that $c+\frac{\nabla\cdot(p\nabla u)}{p}=f$. Then $c=\mathbb{E}_{\vartheta\sim P}[f(\vartheta)]$.

In practice, an approximate solution u gives rise to a control variate $v = \nabla \cdot (p\nabla u)/p$ for use in MCMC.

For more, please see the recent review paper of Anastasiou et al. [2022].

Thank you for your attention!

Collaborators:

Alessandro Barp, François-Xavier Briol, Lawrence Carin, Wilson Chen, Nicolas Chopin, Jon Cockayne, Chris Drovandi, Andrew Duncan, Martin Ehler, Matthew Fisher, Mark Girolami, Jackson Gorham, Manuel Graef, Matt Graham, Toni Karvonen, Jeremias Knoblauch, Lester Mackey, Takuo Matsubara, Antonietta Mira, Chris Nemeth, Steve Niederer, Tui Nolan, Emilio Porcu, Dennis Prangle, Marina Riabiz, Simo Särkkä, Shijing Shi, Leah South, Pawel Swietach, Onur Teymur

References:

- A. Anastasiou et al. Stein's method meets statistics: A review of some recent developments. Statistical Science, 2022. To appear.
- M. Arbel, A. Korba, A. Salim, and A. Gretton. Maximum mean discrepancy gradient flow. NeurIPS, 2019.
- R. Assaraf and M. Caffarel. Zero-variance principle for Monte Carlo algorithms. Physical Review Letters, 83(23):4682, 1999.
- A. Barp, CJO, E. Porcu, and M. Girolami. A Riemann-Stein kernel method. Bernoulli, 2022. To appear.
- D. Belomestny, L. Iosipoi, and N. Zhivotovskiy. Variance reduction via empirical variance minimization: convergence and complexity. arXiv:1712.04667, 2017.
- CJO, M. Girolami, and N. Chopin. Control functionals for Monte Carlo integration. JRSSB, 79(3):695-718, 2017.
- G. Detommaso, T. Cui, Y. Marzouk, R. Scheichl, and A. Spantini. A Stein variational Newton method. In NeurIPS, 2018.
- M. Ehler, M. Gräf, and C. J. Oates. Optimal Monte Carlo integration on closed manifolds. *Statistics and Computing*, 29(6): 1203–1214, 2019.
- M. A. Fisher, T. Nolan, M. M. Graham, D. Prangle, and CJO. Measure transport with kernel Stein discrepancy. AISTATS, 2021.
- W. Gong, Y. Li, and J. M. Hernández-Lobato. Sliced kernelized Stein discrepancy. In ICLR, 2020.
- J. Gorham and L. Mackey. Measuring sample quality with Stein's method. In NeurIPS, 2015.
- J. Gorham and L. Mackey. Measuring sample quality with kernels. In ICML, 2017.
- J. Gorham, A. B. Duncan, S. J. Vollmer, and L. Mackey. Measuring sample quality with diffusions. *AoAP*, 29(5):2884–2928, 2019.



- J. Gorham, A. Raj, and L. Mackey. Stochastic Stein discrepancies. In NeurIPS, 2020.
- M. Gräf, D. Potts, and G. Steidl. Quadrature errors, discrepancies, and their relations to halftoning on the torus and the sphere. SIAM Journal on Scientific Computing, 34(5):A2760–A2791, 2012.
- W. Grathwohl, K.-C. Wang, J.-H. Jacobsen, D. Duvenaud, and R. Zemel. Learning the Stein discrepancy for training and evaluating energy-based models without sampling. In *ICML*, 2020.
- L. Hodgkinson, R. Salomone, and F. Roosta. The reproducing Stein kernel approach for post-hoc corrected sampling. arXiv:2001.09266. 2020.
- T. Hu, Z. Chen, H. Sun, J. Bai, M. Ye, and G. Cheng. Stein neural sampler. arXiv preprint arXiv:1810.03545, 2018.
- J. Huggins and L. Mackey. Random feature Stein discrepancies. In NeurIPS, 2018.
- T. Karvonen and S. Särkkä. Fully symmetric kernel quadrature. SIAM Journal on Scientific Computing, 40(2):A697-A720, 2018.
- T. Karvonen, S. Särkkä, and CJO. Symmetry exploits for Bayesian cubature methods. Statistics and Computing, 29(6): 1231–1248, 2019.
- A. Korba, P.-C. Aubin-Frankowski, S. Majewski, and P. Ablin. Kernel Stein discrepancy descent. In ICML, 2021.
- H. Le, A. Lewis, K. Bharath, and C. Fallaize. A diffusion approach to Stein's method on Riemannian manifolds. arXiv:2003.11497, 2020.
- C. Liu and J. Zhu. Riemannian Stein variational gradient descent for Bayesian inference. In AAAI Conference on AI, number 1, 2018.
- Q. Liu. Stein Variational Gradient Descent as Gradient Flow. In NeurIPS, 2017.
- Q. Liu and J. D. Lee. Black-box importance sampling. In AISTATS, 2017.
- Q. Liu and D. Wang. Stein variational gradient descent: A general purpose Bayesian inference algorithm. In NeurIPS, 2016.
- A. Mira, R. Solgi, and D. Imparato. Zero variance Markov chain Monte Carlo for Bayesian estimators. Statistics and Computing, 23(5):653–662, 2013.
- R. Ranganath, D. Tran, J. Altosaar, and D. Blei. Operator variational inference. In NeurIPS, volume 29, 2016.
- M. Riabiz, W. Chen, J. Cockayne, P. Swietach, S. A. Niederer, L. Mackey, and CJO. Optimal thinning of MCMC output. JRSSB, 2022.
- C.-J. Simon-Gabriel, A. Barp, and L. Mackey. Metrizing weak convergence with maximum mean discrepancies. arXiv:2006.09268, 2020.



- L. F. South, T. Karvonen, C. Nemeth, M. Girolami, and CJO. Semi-exact control functionals from Sard's method. *Biometrika*, 2022.
- B. K. Sriperumbudur, K. Fukumizu, and G. R. Lanckriet. Universality, characteristic kernels and RKHS embedding of measures. Journal of Machine Learning Research, 12(7), 2011.
- C. Stein. A bound for the error in the normal approximation to the distribution of a sum of dependent random variables. In Proceedings of the sixth Berkeley symposium on mathematical statistics and probability, volume 2: Probability theory, pages 583–602. University of California Press, 1972.
- O. Teymur, J. Gorham, M. Riabiz, and CJO. Optimal quantisation of probability measures using maximum mean discrepancy. In AISTATS, 2021.
- H. Wendland. Error estimates for interpolation by compactly supported radial basis functions of minimal degree. Journal of Approximation Theory, 93(2):258–272, 1998.
- L. K. Wenliang and H. Kanagawa. Blindness of score-based methods to isolated components and mixing proportions. arXiv preprint arXiv:2008.10087, 2020.
- W. Xu and G. Reinert. A Stein goodness-of-test for exponential random graph models. In AISTATS, 2021.