# Gaussian process regression model for distribution inputs

UQSay

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## The regression problem for distribution inputs

We are given *n* input/output couples  $(\mu_i, y_i) \in \mathcal{P}(\mathbb{R}) \times \mathbb{R}$ , and we are looking to associate an output to a new input  $\mu_{n+1}$ .



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Our motivations are twofold: we want to deal with regression problems which inputs are

- 1. probability distributions (ex: aggregated or anonymised data, ...)
- 2. functional objects (spectra, histograms, ...)
  - with the nonnegative values and mass 1 restrictions
  - ... which in turn allow the use of tools such as the Wasserstein distance

- 1. Gaussian Process Regression
- 2. Existence of models Stationary kernels on the Wasserstein space
- 3. Maximum-likelihood model selection Asymptotic results
- 4. Numerical performances

## **Gaussian Process Regression**



We chose a random process  $(Y_x)_{x\in\mathbb{R}}$  and consider

 $\hat{Y}(x) := \mathbb{E}(Y_x | Y_{x_1} = y_1, \cdots, Y_{x_n} = y_n)$ 



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Here we need a random process  $(Y_{\mu})_{\mu\in\mathcal{P}(\mathbb{R})}$  to consider

$$\widehat{Y}(\mu) := \mathbb{E}(Y_{\mu}|Y_{\mu_1} = y_1, \cdots, Y_{\mu_n} = y_n)$$

# Existence of models – Stationary kernels on the Wasserstein space

$$W_2(\mu,\nu) := \left(\inf_{\pi \in \Pi(\mu,\nu)} \int_{\mathbb{R}^2} |x-y|^2 d\pi(x,y)\right)^{1/2},$$

where  $\Pi(\mu, \nu)$  is the set of probability distributions on  $\mathbb{R}^2$  with margins  $\mu$  and  $\nu$ .



We obtain a metric space  $\mathcal{W}_2(\mathbb{R})$ .





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distributions on  $\mathbb{R}^{2}$  with margins  $\mu$  and  $\nu$ .  
 $\pi$  We obtain a metric space  $W_{2}(\mathbb{R})$ .

For  $\mu, \nu \in \mathcal{W}_2(\mathbb{R})$  and  $F_{\mu}^{-1}$ ,  $F_{\nu}^{-1}$  the associated quantile functions,

$$W_{2}(\mu,\nu) = \left(\int_{[0,1]} \left(F_{\mu}^{-1}(u) - F_{\nu}^{-1}(u)\right)^{2} du\right)^{1/2}.$$
 (1)

- This optimal coupling, which is specific to the dimension 1 case, allows the numerical evaluation of Wasserstein distances.
- It is also the main ingredient of the proofs of Theorems 1 and 2.

#### Theorem 1 (Fractional Brownian fields)

For ever  $0 \leq H \leq 1$  and  $\sigma_0 \in W_2(\mathbb{R})$ ,

$$K^{H,\sigma_0}(\mu,\nu) = \frac{1}{2} \left( W_2^{2H}(\sigma_0,\mu) + W_2^{2H}(\sigma_0,\nu) - W_2^{2H}(\mu,\nu) \right)$$
(2)

is a covariance function on  $\mathcal{W}_2(\mathbb{R})$ . Moreover, it is nondegenerated if and only if 0 < H < 1.

- We get a fractional Brownian field indexed by W<sub>2</sub>(ℝ). It is a generalisation of the time-indexed fractional Brownian motion, which inherits many enjoyable properties:
- Statistical auto-similarity, path-regularity and long distance memory that are governed by the *Hurst parameter H*.

#### Theorem 2 (Stationary processes)

For every completely monotone  $F : \mathbb{R}^+ \to \mathbb{R}^+$  and  $0 < H \leq 1$ ,

$$(\mu,\nu)\mapsto F\left(W_2^{2H}(\mu,\nu)\right) \tag{3}$$

is a stationary covariance function on  $\mathcal{W}_2(\mathbb{R})$ .

- Recall that F ∈ C<sup>∞</sup>(ℝ<sup>+</sup>, ℝ<sup>+</sup>) is completely monotone if (-1)<sup>n</sup>F<sup>(n)</sup> is nonnegatively valued for every n ∈ N.
- In particular for every  $\sigma^2, \ell > 0$  and  $0 \le H \le 1$ ,

$$\mathcal{K}_{\sigma^2,\ell,H}(\nu_1,\nu_2) = \sigma^2 \exp\left(-\frac{W_2(\nu_1,\nu_2)^{2H}}{\ell}\right) \tag{M}$$

is a valid covariance.

# Maximum-likelihood model selection – Asymptotic results

## Maximum likelihood estimator

- Consider input distributions  $\mu_1, \ldots, \mu_n \in W_2(\mathbb{R})$  and observations  $Y_{\mu_1}, \ldots, Y_{\mu_n}$
- Let  $y = (Y_{\mu_1}, \dots, Y_{\mu_n})^\top$
- Let  $\{K_{\theta}; \theta \in \Theta\}$  be a set of covariance functions on  $\mathcal{W}_2(\mathbb{R})$  with  $\Theta \subset \mathbb{R}^p$
- Let  $R_{ heta}$  be the n imes n matrix  $[K_{ heta}(\mu_i, \mu_j)]_{1 \le i,j \le n}$
- Then the maximum likelihood estimator is

$$\hat{\theta}_{ML} \in \operatorname*{argmin}_{\theta \in \Theta} \left( \mathsf{log}(\mathsf{det}(R_{\theta})) + y^{\top} R_{\theta}^{-1} y \right)$$

#### Condition 1 (Asymptotic expansion framework)

We consider a triangular array of observation points  $\{\mu_1, ..., \mu_n\} = \{\mu_1^{(n)}, ..., \mu_n^{(n)}\}$  so that for all  $n \in \mathbb{N}$  and  $1 \le i \le n, \mu_i$  has support in [i, i + K] with a fixed  $K < \infty$ .

#### Condition 2 (Parametric stationary model)

The model of covariance functions  $\{K_{\theta}, \theta \in \Theta\}$  satisfies

$$\forall \theta \in \Theta, \ K_{\theta}(\mu, \nu) = F_{\theta}(W_2(\mu, \nu)),$$

with  $F_{\theta} : \mathbb{R}^+ \to \mathbb{R}$  and  $\sup_{\theta \in \Theta} |F_{\theta}(t)| \le \frac{A}{1+|t|^{1+\tau}}$  with a fixed  $A < \infty$ ,  $\tau > 1$ .

#### Condition 3 (Well-specified case)

We have observations  $y_i = Y(\mu_i)$ ,  $i = 1, \dots, n$  of the centered Gaussian Process Y with covariance function  $K_{\theta_0}$  for some  $\theta_0 \in \Theta$ .

#### Condition 4 (Asymptotical nondegeneracy)

The sequence of matrices  $R_{\theta} = (K_{\theta}(\mu_i, \mu_j))_{1 \le i,j \le n}$  satisfies

 $\lambda_{\inf}(R_{\theta}) \geq c$ 

for a fixed c > 0, where  $\lambda_{inf}(R_{\theta})$  denotes the smallest eigenvalue of  $R_{\theta}$ .

### Condition 5 (First sampling condition)

 $\forall \alpha > 0$ ,

$$\liminf_{n\to\infty}\inf_{\|\theta-\theta_0\|\geq\alpha}\frac{1}{n}\sum_{i,j=1}^n\left[\mathcal{K}_\theta(\mu_i,\mu_j)-\mathcal{K}_{\theta_0}(\mu_i,\mu_j)\right]^2>0.$$

#### Theorem 3 (Consistency of MLE)

Under conditions 1 to 5, the maximum-likelihood estimator is consistent, that is to say:

$$\hat{\theta}_{ML} \xrightarrow[n \to \infty]{\mathbb{P}} \theta_0.$$

## Supplementary conditions

#### Condition 6 (Model regularity)

- $\forall t \geq 0, F_{\theta}(t) \text{ is } C^{1} \text{ with respect to } \theta \text{ and verifies}$  $\sup_{\theta \in \Theta} \max_{i=1,\cdots,p} \left| \frac{\partial}{\partial \theta_{i}} F_{\theta}(t) \right| \leq \frac{A}{1+t^{1+\tau}}, \text{ where } A, \tau \text{ are defined in}$ Condition 2.
- For every  $t \ge 0$ ,  $F_{\theta}(t)$  is  $C^3$  with respect to  $\theta$  and  $\forall q \in \{2,3\}$ ,  $\forall i_1 \cdots i_q \in \{1, \cdots p\}$ ,

$$\sup_{\theta\in\Theta}\left|\frac{\partial}{\partial\theta_{i_1}}\cdots\frac{\partial}{\partial\theta_{i_q}}\mathsf{F}_\theta(t)\right|\leq \frac{A}{1+|t|^{1+\tau}}.$$

#### Condition 7 (Second sampling condition)

$$\forall (\lambda_1 \cdots, \lambda_p) \neq (0, \cdots, 0),$$
  
$$\liminf_{n \to \infty} \frac{1}{n} \sum_{i,j=1}^n \left( \sum_{k=1}^p \lambda_k \frac{\partial}{\partial_{\theta_k}} K_{\theta_0}(\mu_i, \mu_j) \right)^2 > 0.$$

14

#### Theorem 4

Let  $M_{ML}$  be the  $p \times p$  matrix defined by

$$(M_{ML})_{i,j} = \frac{1}{2n} \operatorname{Tr} \left( R_{\theta_0}^{-1} \frac{\partial R_{\theta_0}}{\partial \theta_i} R_{\theta_0}^{-1} \frac{\partial R_{\theta_0}}{\partial \theta_j} \right).$$

Under conditions 1 to 7, the maximum-likelihood estimator is asymptotically normal:

$$\sqrt{n} \ M_{ML}^{1/2} \left( \hat{\theta}_{ML} - \theta_0 \right) \xrightarrow[n \to \infty]{\mathcal{L}} \mathcal{N}(0, I_{\rho}).$$

Moreover

$$0 < \liminf_{n \to \infty} \lambda_{\min}(M_{ML}) \leq \limsup_{n \to \infty} \lambda_{\max}(M_{ML}) < +\infty.$$

#### **Proposition 1**

Assume that Conditions 2 and 6 hold, that for  $\theta \neq \theta_0$ ,  $F_{\theta}$  and  $F_{\theta_0}$  are not equal everywhere on  $\mathbb{R}^+$ , and that there does not exist  $(\lambda_1, ..., \lambda_p) \neq (0, ..., 0)$  so that  $\sum_{i=1}^p \lambda_i (\partial/\partial \theta_i) F_{\theta_0}$  is the zero function on  $\mathbb{R}^+$ .

Let  $(Z_i)_{i \in \mathbb{Z}}$  be iid, centred Gaussian processes on  $\mathbb{R}$  with continuous trajectories, and stationary covariance  $C_0(u - v)$ .

Assume that  $\hat{C}_0(w)|w|^{2p}$  is bounded away from 0 and  $\infty$  as  $|w| \to \infty$ . Let K > 1 be fixed. For  $i \in \mathbb{N}$ , let  $\mu_i$  be the measure with density

$$f_i(t) = \frac{e^{Z_i(t-i)}}{\int_i^{i+K} e^{Z_i(t-i)dt}} \mathbb{1}_{[i,i+K]}(t).$$

Then, almost surely, with the sequence of random probability measures  $\{\mu_1, ..., \mu_n\}$ , Conditions 5 and 7 hold.

#### Theorem 5

Under conditions 1 to 7, the Kriging estimator under the ML-estimated parameter  $\hat{\theta}_{ML}$  is asymptotically optimal:

$$orall \mu \in \mathcal{W}_2(\mathbb{R}), \; \left| \hat{Y}_{\hat{ heta}_{ML}}(\mu) - \hat{Y}_{ heta_0}(\mu) 
ight| = o_{\mathbb{P}}(1).$$

## Numerical performances

• Denote by  $m_k(\nu)$  the order k moment of  $\nu$ . We consider

$$F: \mathcal{W}_{2}(\mathbb{R}) \to \mathbb{R}$$
$$F(\nu) = \frac{m_{1}(\nu)}{0.05 + \sqrt{m_{2}(\nu) - m_{1}(\nu)^{2}}},$$
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which we are going to regress.

- Let us generate normal random variables ν<sub>1</sub>, · · · , ν<sub>100</sub>, with means and variances drawn uniformly at random, randomly perturbed to exhibit irregularities.
- We estimate  $\hat{\sigma}^2, \hat{\ell}, \hat{H}$  by maximising the maximum likelihood for the parametric model:

$$\mathcal{K}_{\sigma^2,\ell,H}(\nu_1,\nu_2) = \sigma^2 \exp\left(-\frac{W_2(\nu_1,\nu_2)^{2H}}{\ell}\right).$$
 (5)

 We evaluate the method on a test dataset (ν<sub>t,i</sub>)<sup>500</sup><sub>i=1</sub> which is generated in a same way as the ν<sub>i</sub>,

• We evaluate the method on a test dataset  $(\nu_{t,i})_{i=1}^{500}$  which is generated in a same way as the  $\nu_i$ , with the criteria:

$$RMSE^{2} = \frac{1}{500} \sum_{i=1}^{500} \left( F(\nu_{t,i}) - \hat{F}(\nu_{t,i}) \right)^{2},$$
$$CIR_{\alpha} = \frac{1}{500} \sum_{i=1}^{500} \mathbf{1} \left\{ \left| F(\nu_{t,i}) - \hat{F}(\nu_{t,i}) \right| \le q_{\alpha} \hat{\sigma}(\nu_{t,i}) \right\}.$$

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modèle	RMSE	$CIR_{0.9}$
"Wasserstein"	0.094	0.92
"Legendre" ordre 5	0.49	0.92
"Legendre" ordre 10	0.34	0.89
"Legendre" ordre 15	0.29	0.91
"PCA" ordre 5	0.63	0.82
"PCA" ordre 10	0.52	0.87
"PCA" ordre 15	0.47	0.93

## Two stage sampling

- In [6], Poczos and al. try to predict outputs *S*(*P*) corresponding to distributions *P*, where only samples of those are available.
  - 1. Starting with a kernel smoothing  $\hat{P}, \hat{P}_1, \cdots, \hat{P}_n$  of the empirical distributions.
  - Then the prediction \$\hfrac{S}(\hfrac{P})\$ of \$S(P)\$ is obtained by a weighted average of \$S(P\_1), ..., \$S(P\_n)\$. Weights are obtained by applying some kernel to the \$L^1\$ distance between densities \$\hfrac{P}{p}\$ and \$\hfrac{P\_1}{p\_1}, \dots, \$\hfrac{P\_n}{p\_n}\$.

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- We add a nugget term to our covariance to accomodate for the differenve between S(P) and S(P̂):

$$K_{\sigma^{2},\ell,H,\delta}(\nu_{1},\nu_{2}) = \sigma^{2} \exp\left(-\frac{W_{2}(\nu_{1},\nu_{2})^{2H}}{\ell}\right) + \delta \mathbf{1}_{\{W_{2}(\nu_{1},\nu_{2})=\mathbf{0}\}},$$

and obtain the following results:

model	RMSE	CIR <sub>0.9</sub>
"distribution"	0.21	0.91
"kernel regression"	0.93	

- For distributions on ℝ<sup>d</sup>, d > 1, the nice representation with the quantile functions does not hold any longer
- Our kernels do not remain non-negative definite when extended to distributions on  $\mathbb{R}^d$ , d>1
- We suggest to compute *L*<sup>2</sup> norms between optimal transport maps to obtain non-negative definite kernels.

## F. Bachoc.

Asymptotic analysis of the role of spatial sampling for covariance parameter estimation of Gaussian processes. *Journal of Multivariate Analysis*, 125:1–35, 2014.

F. Bachoc, F. Gamboa, J.-M. Loubes, and N. Venet.
 Gaussian process regression model for distribution inputs.
 IEEE Transactions on Information Theory, 64(10):6620–6637, 2018.

F. Bachoc, A. Suvorikova, D. Ginsbourger, J.-M. Loubes, and V. Spokoiny.

Gaussian processes for multidimensional distribution inputs via optimal transport and Hilbertian embedding.

Electronic Journal of Statistics, 14(2):2742-2772, 2020.



C. Berg, J. P. R. Christensen, and P. Ressel. Harmonic analysis on semigroups. Springer-Verlag, 1984.



J. Istas.

Manifold indexed fractional fields.

ESAIM Probab. Stat., 16:222–276, 2012.

B. Póczos, A. Singh, A. Rinaldo, and L. Wasserman. **Distribution-free distribution regression.** 

In In Proceedings of the 16th International Conference on Artificial Intelligence and Statistics, volume 31 of JMLR Proceedings, pages 507–515, 2013.

## References iii

N. Venet.

Nonexistence of fractional brownian fields indexed by cylinders.

#### arXiv preprint, 2016.



N. Venet.

On the existence of fractional brownian fields indexed by manifolds with closed geodesics.

arXiv preprint, 2016.



C. Villani.

Optimal transport: old and new, volume 338.

Springer Science & Business Media, 2009.

## Thank you for your attention