# Machine learning - Master ICFP 2019-2020 <br> Kernel methods 

Francis Bach

March 6, 2020

To learn more about the topic of this lecture, please look at the following documents:

- http://cbio.ensmp.fr/~jvert/svn/kernelcourse/slides/master/master.pdf
- http://www.gatsby.ucl.ac.uk/~gretton/coursefiles/lecture4_introToRKHS.pdf
- http://www.di.ens.fr/~fbach/rasma_fbach.pdf

In this course, we often focused on prediction methods which are linear, that is, the input data are vectors (i.e., $x \in \mathbb{R}^{d}$ ) and the prediction function is linear: $f(x)=w^{\top} x$ for $w \in \mathbb{R}^{d}$. In this situation, given data $\left(x_{i}, y_{i}\right), i=1, \ldots, n$, the vector $w$ is obtained by minimizing

$$
\frac{1}{n} \sum_{i=1}^{n} \ell\left(y_{i}, w^{\top} x_{i}\right)+\lambda \Omega(w) .
$$

Classical examples are logistic regression or least-squares regression.
These methods look at first sight of limited practical significance, because:

- Input data may not be vectors.
- Relevant prediction functions may not be linear.

The goal of kernel methods is to go beyond theses limitations while keeping the good aspects. The underlying principle is to replace $x$ by any function $\varphi(x) \in \mathbb{R}^{d}$, explicitly ou implicitly, and consider linear predictions in $\Phi(x)$, i.e., $f(x)=w^{\top} \varphi(x)$. We call $\varphi(x)$ the "feature" associated to $x$.

Example. Polynomial regression of degree $r$, by considering $x \in \mathbb{R}^{d}$ and

$$
\varphi(x)=\left(x_{1}^{\alpha_{1}} \cdots x_{d}^{\alpha_{d}}\right)_{\sum_{i=1}^{d} \alpha_{i}=r} .
$$

In this situation, $p=\binom{d+r-1}{r}$ (number of $k$-combinations with repetitions from a set with cardinality $d$ ), can be too big for an explicit representation to be feasible.

WARNING. The type of kernel is different from the ones in lecture 2. The ones here are "positive definite"; the ones from lecture 2 are "non-negative". See more details in https://francisbach.com/ cursed-kernels/

## 1 Representer theorem

Theorem 1 (Representer theorem, 1971).
Let $\varphi: X \rightarrow \mathbb{R}^{d}$. Let $\left(x_{1}, \ldots, x_{n}\right) \in X^{n}$, and assume $\Psi: \mathbb{R}^{n+1} \rightarrow \mathbb{R}$ strictly increasing with respect to the last variable, then the minimum of $\Psi\left(w^{\top} \varphi\left(x_{1}\right), \ldots, w^{\top} \varphi\left(x_{n}\right), w^{\top} w\right)$ is attained for $w=\sum_{i=1}^{n} \alpha_{i} \Phi\left(x_{i}\right)$ with $\alpha \in \mathbb{R}^{n}$.

Proof Let $w \in \mathbb{R}^{d}$, and $\mathcal{F}_{D}=\left\{\sum \alpha_{i} \Phi\left(x_{i}\right) / \alpha \in \mathbb{R}^{n}\right\}$. Let $w_{D} \in \mathcal{F}_{D}$ and $w_{\perp} \in \mathcal{F}_{D}^{\perp}$ such that $w=w_{D}+w_{\perp}$, then $\forall i, w^{\top} \varphi\left(x_{i}\right)=w_{D}^{\top} \varphi\left(x_{i}\right)+w_{\perp}^{\top} \varphi\left(x_{i}\right)$ with $w_{\perp}^{\top} \varphi\left(x_{i}\right)=0$.
From Pythagoreas theorem, we get: $w^{\top} w=w_{D}^{\top} w_{D}^{2}+w_{\perp}^{\top} w_{\perp}$. Therefore we have:

$$
\begin{aligned}
\Psi\left(w^{\top} \varphi\left(x_{1}\right), \ldots, w^{\top} \varphi\left(x_{n}\right), w^{\top} w\right) & =\Psi\left(w_{D}^{\top} \varphi\left(x_{1}\right), \ldots, w_{D}^{\top} \varphi\left(x_{n}\right), w_{D}^{\top} w_{D}+w_{\perp}^{\top} w_{\perp}\right) \\
& \geq \Psi\left(w_{D}^{\top} \varphi\left(x_{1}\right), \ldots, w_{D}^{\top} \varphi\left(x_{n}\right), w_{D}^{\top} w_{D}\right) .
\end{aligned}
$$

Thus

$$
\inf _{w \in \mathbb{R}^{d}} \Psi\left(w^{\top} \varphi\left(x_{1}\right), \ldots, w^{\top} \varphi\left(x_{n}\right), w^{\top} w\right)=\inf _{w \in \mathcal{F}_{D}} \Psi\left(w^{\top} \varphi\left(x_{1}\right), \ldots, w^{\top} \varphi\left(x_{n}\right), w^{\top} w\right)
$$

Corollary 1 For $\lambda>0, \min _{w \in \mathbb{R}^{d}} \frac{1}{n} \sum \ell\left(y_{i}, w^{\top} \varphi\left(x_{i}\right)\right)+\frac{\lambda}{2} w^{\top} w$ is attained at $w=\sum_{i=1}^{n} \alpha_{i} \varphi\left(x_{i}\right)$.

- It is important to note that there is no assumption on $\ell$ (no convexity).
- This result is extendable to Hilbert spaces (RKHS).
- We have: $\forall j \in\{1, \ldots, n\}, w^{\top} \varphi\left(x_{j}\right)=\sum_{i=1}^{n} \alpha_{i} k\left(x_{i}, x_{j}\right)=(K \alpha)_{j}$ where $K$ is the kernel matrix and $w^{\top} w=\alpha^{\top} K \alpha$. We can then write:

$$
\min _{w \in \mathbb{R}^{d}} \frac{1}{n} \sum \ell\left(y_{i}, w^{\top} \varphi\left(x_{i}\right)\right)+\frac{\lambda}{2} w^{\top} w=\min _{\alpha \in \mathbb{R}^{n}} \frac{1}{n} \sum \ell\left(y_{i},(K \alpha)_{i}\right)+\frac{\lambda}{2} \alpha^{\top} K \alpha .
$$

For a test point, we have $f(x)=\sum_{i=1}^{n} \alpha_{i} k\left(x, x_{i}\right)$.
The kernel trick allows to:

- replace $\mathbb{R}^{d}$ by $\mathbb{R}^{n}$; this is interesting when $d$ is very large.
- separate the representation problem (design a kernel on a set $X$ ) and algorithms and analysis (which only use the kernel matrix $K$ ).


## 2 Kernels

- Definition: $k$ is a positive definite kernel if and only if all kernel matrices are positive semi-definite.

Theorem 2 (Aronszajn, 1950)
$k$ is a positive definite kernel if and only if there exists a Hilbert space $\mathcal{F}$, and $\Phi: \mathcal{X} \rightarrow \mathcal{F}$ such that $\forall x, y, k(x, y)=\langle\Phi(x), \Phi(y)\rangle$.

- $\mathcal{F}$ is called the "feature space", and $\varphi$ the "feature map".
- Simple properties (to be done as exercises): the sum and product of kernels are kernels. What are their associated feature space and feature map?
- Linear kernel: $k(x, y)=x^{\top} y$
- Polynomial kernel: the kernel $k(x, y)=\left(x^{\top} y\right)^{r}$ can be expanded as:

$$
k(x, y)=\left(\sum_{i=1}^{d} x_{i} y_{i}\right)^{r}=\sum_{\alpha_{1}+\ldots+\alpha_{p}=r}\binom{r}{\alpha_{1}, \ldots, \alpha_{p}} \underbrace{\left(x_{1} y_{1}\right)^{\alpha_{1}} \ldots\left(x_{p} y_{p}\right)^{\alpha_{p}}}_{\left(x_{1}^{\alpha_{1}} \ldots x_{p}^{\alpha_{p}}\right)\left(y_{1}^{\alpha_{1}} \ldots y_{p}^{\alpha_{p}}\right)}
$$

We have: $\Phi(x)=\left\{\binom{r}{\alpha_{1}, \ldots, \alpha_{p}}^{\frac{1}{2}} x_{1}^{\alpha_{1}} \ldots x_{p}^{\alpha_{p}}\right\}$. Exercise: how can we go beyond homogeneous polynomials?

- Translation-invariant kernels on $[0,1] . k(x, y)=q(x-y)$ where $q$ is 1-periodic. $k$ is a positive definite kernel if and only if the Fourier series of $q$ is non-negative (using the complex representation), i.e.,

$$
k(x, y)=\nu_{0}+\sum_{m \geqslant 1} 2 \nu_{m} \cos 2 \pi m x \cos 2 \pi m y+2 \nu_{m} \sin 2 \pi m x \sin 2 \pi m y
$$

with $\nu \geqslant 0$.
The (infinite-dimensional) feature vector is composed of $\nu_{0}^{1 / 2}$, and of $\sqrt{2 \nu_{m}} \cos 2 \pi m x$ and $\sqrt{2 \nu_{m}} \sin 2 \pi m x$, for $m \geqslant 1$.
If $f(x)$ can be written $f(x)=\Phi(x)^{\top} w$, then

$$
\|w\|^{2}=\left(\int_{0}^{1} f(x)\right)^{2}+\sum_{m \geqslant 1} \frac{2}{\nu_{m}}\left(\int_{0}^{1} f(x) \cos 2 \pi m x\right)^{2}+\frac{2}{\nu_{m}}\left(\int_{0}^{1} f(x) \sin 2 \pi m x\right)^{2} .
$$

For $\nu_{m}=\frac{1}{m^{2 s}}, m \geqslant 1$, this norm is equal to

$$
\|w\|^{2}=\left(\int_{0}^{1} f(x)\right)^{2}+\frac{1}{(2 \pi)^{2 s}} \int_{0}^{2}\left|f^{(s)}(x)\right|^{2} d x
$$

and the kernel has an analytical expression $k(x, y)=\nu_{0}+(-1)^{s-1} \frac{(2 \pi)^{2 s}}{(2 s)!} B_{2 s}(\{x-y\})$, where $B_{2 s}$ is Bernoulli's polynomial.

- Translation-invariant kernels on $\mathbb{R}^{d}: X=\mathbb{R}^{d}, k(x, y)=q(x-y)$ with $q: \mathbb{R}^{d} \rightarrow \mathbb{R}$.

Theorem 3 (Böchner): $k$ is positive definite $\Leftrightarrow q$ is the Fourier transform of a non-negative Borel measure $\Leftarrow q \in L^{1}$ and its Fourier transform is non-negative.

Proof (partial) Let $x_{1}, \ldots x_{n} \in \mathbb{R}^{d}$, let $\alpha_{1}, . ., \alpha_{n} \in \mathbb{R}$,

$$
\begin{aligned}
\sum \alpha_{s} \alpha_{j} k\left(x_{s}, x_{j}\right) & =\sum \alpha_{s} \alpha_{j} q\left(x_{s}-x_{j}\right) \\
& =\sum \alpha_{s} \alpha_{j} \int \exp ^{-i \omega^{\top}\left(x_{s}-x_{j}\right)} d \mu(\omega) \\
& =\int\left(\sum \alpha_{s} \alpha_{j} \exp ^{-i \omega^{\top} x_{s}} \overline{\exp ^{-i \omega^{\top} x_{j}}}\right) d \mu(\omega) \\
& =\int\left|\sum \alpha_{s} \exp ^{-i \omega^{\top} x_{s}}\right|^{2} d \mu(\omega) \geq 0 .
\end{aligned}
$$

Consruction of the norm. Intuitive (non-rigorous) reasoning: if $q$ is in $L^{1}$, then $\widehat{q}(\omega)$ exists and, with $d \mu(\omega)=\widehat{q}(\omega) d \omega$, we have an explicit representation of

$$
k(x, y)=\int\left\langle\sqrt{\widehat{q}(\omega)} \exp ^{-i \omega^{\top} x}, \sqrt{\widehat{q}(\omega)} \exp ^{-i \omega^{\top} x}\right\rangle d \omega=\int\left\langle\varphi_{\omega}(x), \varphi_{\omega}(y)\right\rangle d \omega=\langle\varphi(x), \varphi(y)\rangle .
$$

If we consider $f(x)=\int \varphi_{\omega}(x) w_{\omega} d \omega$, then $w_{\omega}=\widehat{f}(\omega) / \sqrt{\widehat{q}(\omega)}$, and the squared norm of $w$ is equal to $\int \frac{|\widehat{f}(w)|^{2}}{\widehat{q}(w)} d w$, where $\widehat{f}$ denotes the Fourier transform of $f$.
Examples: Exponential kernel $\exp (-\alpha|x-y|)$ and Gaussian kernel $\exp \left(-\alpha|x-y|^{2}\right)$.

- Many applications of the kernel trick!
- Exercise: show that on $X=\mathbb{R}^{+}, k(x, y)=\min (x, y)$ and $k(x, y)=\frac{x y}{x+y}$ are positive definite kernels.
- Non vectorial data (sequences, graphes, images).
- Exercise: for $X$ the set of all subsets of a given set $V$, show that $k(A, B)=\frac{|A \cap B|}{|A \cup B|}$ is a positive definite kernel.
- Examples of kernels on sequences


## 3 Ridge regression (mostly as an exercise)

We consider the optimization problem:

$$
\min _{w \in \mathbb{R}^{d}} \frac{1}{2 n}\|y-\Phi w\|_{2}^{2}+\frac{\lambda}{2}\|w\|_{2}^{2} .
$$

We can solve it in two ways (done as an execise):

1. Direct : $\min _{w \in \mathbb{R}^{d}} \frac{1}{2 n}\|y-\Phi w\|_{2}^{2}+\frac{\lambda}{2}\|w\|_{2}^{2}$
2. With representer theorem : $\min _{\alpha \in \mathbb{R}^{n}} \frac{1}{2 n}\|y-K \alpha\|_{2}^{2}+\frac{\lambda}{2} \alpha^{\top} K \alpha$
3. Using the representer theorem:
gradient with respect to $\alpha: \frac{1}{n} K(K \alpha-y)+\lambda K \alpha=0 \Leftrightarrow\left(K^{2}+n \lambda K\right) \alpha=K y \Leftrightarrow K((K+n \lambda I) \alpha-y)=0$. If $K$ is non invertible, the solution is not unique : $\alpha=(K+n \lambda I)^{-1} y+\operatorname{Ker}(K)$. However the prediction is unique : $K \alpha=K(K+n \lambda I)^{-1} y$.
4. Direct method: minimizing with respect to $w$
gradient w.r.t. $w: \frac{1}{n} \Phi^{\top}(\Phi w-y)$
This leads to $w=\left(\frac{1}{n} \Phi^{\top} \Phi+\lambda I\right)^{-1} \frac{1}{n} \Phi^{\top} y \Leftrightarrow \Phi f=\Phi\left(\frac{1}{n} \Phi^{\top} \Phi+\lambda I\right)^{-1} \frac{1}{n} \Phi^{\top} y$.
With $K=\Phi \Phi^{\top}$, we get :

$$
\overbrace{\Phi \Phi^{\top}(\underbrace{\Phi \Phi^{\top}}_{n \times n}+n \lambda I)^{-1} y}^{\text {kernel }}=\overbrace{\Phi(\underbrace{\Phi^{\top} \Phi}_{d \times d}+n \lambda I)^{-1} \Phi^{\top} y}^{\text {direc }}
$$

This is simply:
Lemma 1 (matrix inversion lemma) $\forall A$ matrix, $\left(A A^{\top}+I\right)^{-1} A=A\left(A^{\top} A+I\right)^{-1}$
There is thus an "equivalence" betweeen this lemma and the representer theorem.

## 4 Complexity of linear algebra computations

If $K \in \mathbb{R}^{n \times n}$ and $L \in \mathbb{R}^{n \times n}$ are two matrices

- computing $K L$ has complexity $O\left(n^{3}\right)$
- computing $K^{-1}$ has complexity $O\left(n^{3}\right)$
- computing $K y$ has complexity $O\left(n^{2}\right)$
- Solving $K^{-1} y$ has complexity $O\left(n^{3}\right)$
- Decomposing $K$ in eigenvalues / eigenvectors $O\left(n^{3}\right)$
- Largest eigenvector: $O\left(n^{2}\right)$

Low-rank approximation

- Eigenvector basis (complexity $O\left(n^{2} r\right)$ )
- Orthogonal projection on first $r$ columns: $O\left(n r^{2}\right)$


## 5 Using distances in feature space (exercise)

Given sets of negative examples $x_{i}, i \in I_{-}$and positive examples $x_{i}, i \in I_{+}$, we consider the average $\mu_{-}$of negative points and the average $\mu_{+}$of positive points in the feature space.
(1) For a testing point $x$, compute $\left\|\Phi(x)-\mu_{+}\right\|^{2}$ and $\left\|\Phi(x)-\mu_{-}\right\|^{2}$.
(2) We classify $x$ as positive if $\left\|\Phi(x)-\mu_{+}\right\|^{2}>\left\|\Phi(x)-\mu_{-}\right\|^{2}$. Relate the classification rule to existing classifiers. We can consider that $\left|I_{+}\right|=\left|I_{-}\right|$.

