${\it Modern Statistical Methods-Summary}$

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Classical models rely on so-called "large n asymptotics" (where n is the sample size). This course focuses on the scenario where p, the number of variables, is larger or about as large as n. In this case, the classical theory breaks down, so we need new methods.

1 Kernel machines

We represent data are pairs $(Y_i, x_i) \in \mathbb{R} \times \mathbb{R}^p$ (i = 1, ..., n). The random variables Y_i are called the responses, and the (fixed) variables x_i are called predictors.

Recap 1.1. Let $X = (X_1, \dots, X_n)^{\top}$ be a multivariate random variable. Its distribution function is given by

$$F_X \colon \mathbb{R}^n \to [0,1] \colon \mathbf{x} \mapsto \mathbb{P}(X_1 \le x_1, \dots, X_n \le x_n).$$

Its expected value is given by

$$\mathbb{E}[X] := (\mathbb{E}[X_1], \dots, \mathbb{E}[X_n])^{\top} \in \mathbb{R}^n.$$

Its covariance matrix is given by

$$\operatorname{Var}[X] = \mathbb{E}[(X - \mathbb{E}[X])(X - \mathbb{E}[X])^{\top}] = \mathbb{E}[XX^{\top}] - \mathbb{E}[X]\mathbb{E}[X]^{\top}.$$

The matrix Var[X] is symmetric positive semidefinite and satisfies $(Var[X])_{ij} = Cov(X_i, X_j)$.

Definition 1.2. In a *linear model*, we assume that

$$Y_i = x_i^{\mathsf{T}} \beta^0 + \varepsilon_i \quad (i = 1, \dots, n).$$

where $\beta \in \mathbb{R}^p$ is unknown and the multivariate random variable $\varepsilon = (\varepsilon_1, \dots, \varepsilon_n)$ satisfies $\mathbb{E}(\varepsilon) = 0$ and $\operatorname{Var}(\varepsilon) = \sigma^2 I$.

Definition 1.3. For an estimator $\tilde{\beta}$ of β^0 , its mean squared error (MSE) is given by

$$\mathbb{E}_{\beta^0,\sigma^2} \Big[(\tilde{\beta} - \beta^0) (\tilde{\beta} - \beta^0)^\top \Big] = \operatorname{Var}(\tilde{\beta}) + \Big[\mathbb{E}(\tilde{\beta} - \beta^0) \Big] \Big[\mathbb{E}(\tilde{\beta} - \beta^0) \Big]^\top.$$

Note that if $\tilde{\beta}$ is unbiased, the second term will disappear and the MSE is simply the variance.

Recap 1.4. The maximum likelihood estimator (MLE) in this model is the ordinary least squares (OLS) estimator $\hat{\beta}^{\text{OLS}} = (X^{\top}X)^{-1}X^{\top}Y$, where the design matrix $X \in \mathbb{R}^{n \times p}$ is the matrix whose rows are the vectors x_i . This estimator only exists if X has full column rank, so in particular, it is required that p < n.

The Cramér-Rao lower bown states that, out of all unbiased estimators, the MLE has the optimal variance asymptotically (i.e., for $n \to \infty$).

1.1 Ridge regression

Definition 1.5. Let $\lambda \geq 0$, and let $\mathbf{1} \in \mathbb{R}^n$ be the all-ones vector. Then we define the *ridge regression* estimators

$$(\hat{\mu}_{\lambda}^{\mathrm{R}}, \hat{\beta}_{\lambda}^{\mathrm{R}}) \coloneqq \underset{(\mu,\beta) \in \mathbb{R} \times \mathbb{R}^p}{\arg \min} \Big\{ \big\| Y - \mu \mathbf{1} - X\beta \big\|^2 + \lambda \big\| \beta \big\|^2 \Big\},$$

where the used norm is the 2-norm. The parameter λ is called the regularisation parameter.

The parameter λ represents a penalty for large coefficients in the design matrix. The intercept is not penalised — this is because a shift in units should not affect the fitted values. However, $X\hat{\beta}$ is not invariant under scale transformations, so it is common practice to centre the columns of X to have mean 0, and then scale them to have ℓ_2 -norm \sqrt{n} .

After that, we can compute $\hat{\mu}_{\lambda}^{R}$ by taking the derivative:

$$\|Y - \mu \mathbf{1} - X\beta\|^2 = \sum_{i} (Y_i - \mu - \sum_{j} X_{ij}\beta_j)^2.$$
$$\frac{\partial}{\partial \mu} \|Y - \mu \mathbf{1} - X\beta\|^2 = -2\sum_{i} \left(Y_i - \mu - \sum_{j} X_{ij}\beta_j \right).$$

Setting this derivative equal to 0 yields

$$-2\sum_{i} \left(Y_{i} - \mu - \sum_{j} X_{ij} \beta_{j} \right) = 0$$

$$\sum_{i} Y_{i} - n\mu - \sum_{j} \beta_{j} \left(\sum_{i} X_{ij} \right) = 0$$

$$\sum_{i} Y_{i} - n\mu = 0$$

$$\mu = \frac{1}{n} \sum_{i} Y_{i} = \bar{Y}.$$

Therefore we conclude $\hat{\mu}_{\lambda}^{R} = \bar{Y}$. After centering the responses (i.e. replacing Y_i by $Y_i - \bar{Y}$), the problem can be reduced to

$$\hat{\beta}_{\lambda}^{\mathrm{R}} = \operatorname*{arg\,min}_{\beta \in \mathbb{R}^{p}} Q(\beta) \coloneqq \operatorname*{arg\,min}_{\beta \in \mathbb{R}^{p}} \Big[\|Y - X\beta\|^{2} + \lambda \|\beta\|^{2} \Big].$$

Since $Q(\beta)$ is convex quadratic, there is a unique root, and to find it we compute

$$\nabla_{\beta} Q(\beta) = 2X^{\top} (Y - X\beta) + 2\lambda \beta = 0 \iff \beta = (X^{\top} X + \lambda I)^{-1} X^{\top} Y.$$

We conclude that $\hat{\beta}_{\lambda}^{\mathrm{R}} = (X^{\top}X + \lambda I)^{-1}X^{\top}Y$. Note that, even if X does not have full column rank, this estimator exists for all $\lambda > 0$. In fact, for λ sufficiently small, the ridge estimator outperforms the MLE in terms of mean squared error:

Theorem 1.6. Fix β^0 , σ^2 , and assume that $\hat{\beta}^{OLS}$ exists (i.e., X has full column rank). For some $\lambda > 0$ sufficiently small, it holds that the MSE of $\hat{\beta}^{OLS}$ minus the MSE of $\hat{\beta}^{R}_{\lambda}$ is positive definite.

Proof. This is simply writing out the MSE's. In the end, we find that the result holds for $0 < \lambda < 2\sigma^2/\|\beta^0\|^2$.

1.1.1 The SVD and PCA

Recap 1.7. Recall that any $X \in \mathbb{R}^{n \times p}$ can be factorised as $X = UDV^{\top}$, where U, V are $n \times n$ and $p \times p$ orthogonal matrices respectively, and $D \in \mathbb{R}^{n \times p}$ satisfies $D_{11} \ge \cdots \ge D_{mm} \ge 0$ where $m := \min(n, p)$, and all other entries of D are 0. This is called the *singular value decomposition* or SVD of X.

If n > p, we can replace U by its first p columns and D by its first p rows to produce the so-called thin SVD of X. Then $U \in \mathbb{R}^{n \times p}$ has orthogonal columns (so $U^{\top}U = I$) and $D \in \mathbb{R}^{p \times p}$

Suppose $n \geq p$ and let $X = UDV^{\top}$ be the thin SVD of our design matrix X. Then we can write the fitted values from the ridge regression as follows:

$$\begin{split} X\hat{\beta}_{\lambda}^{\mathrm{R}} &= X(X^{\top}X + \lambda I)^{-1}X^{\top}Y \\ &= UDV^{\top}(VD^{2}V^{\top} + \lambda I)^{-1}VDU^{\top}Y \\ &= UDV^{\top}\big(V(D^{2} + \lambda I)V^{\top}\big)^{-1}VDU^{\top}Y \\ &= UD(D^{2} + \lambda I)^{-1}DU^{\top}Y \\ &= UD^{2}(D^{2} + \lambda I)^{-1}U^{\top}Y \\ &= \sum_{j=1}^{p} \frac{D_{jj}^{2}}{D_{jj}^{2} + \lambda}U_{j}U_{j}^{\top}Y. \end{split}$$

Note that for OLS ($\lambda = 0$), this is simply the projection of Y onto the column space of X (if X has full column rank). If $\lambda > 0$, Y is still projected onto the column space of X, but the projection is shrunk in the directions of the left singular vectors, and the lower the corresponding singular value, the higher the shrinkage.

Principal component analysis Consider $v \in \mathbb{R}^p$ with norm 1, then since the columns of X have been centered, the sample mean of Xv is 0, and the sample variance is therefore

$$\frac{1}{n}\sum_i (Xv)_i^2 = \frac{1}{n}(Xv)^\top Xv = \frac{1}{n}v^\top X^\top Xv = \frac{1}{n}v^\top VD^2 V^\top v.$$

Writing $a = V^{\top}v$ (with ||a|| = 1), we find

$$\frac{1}{n}v^{\top}VD^{2}V^{\top}v = \frac{1}{n}a^{\top}D^{2}a = \frac{1}{n}\sum_{j}a_{j}^{2}D_{jj}^{2}$$

Therefore, we see that the above is maximised if $a = \pm e_1$, or equivalently $v = \pm V_1$. Therefore, V_1 determines which combination of columns of X has the largest variance (subject to having norm 1), and $XV_1 = D_{11}U_1$ is known as the first principal component of X. Analogously, it can be shown that $D_{22}U_2, \ldots, D_{pp}U_p$ have maximum variance D_{jj}^2/n , subject to being orthonormal to all earlier principal components.

We see that ridge regression shrinks Y most in the smaller principal components of X. Therefore it will work well if most of the information is in the larger principal components of X.

A comment on computation By analogous calculations as before, one cam compute $\hat{\beta}_{\lambda}^{R} = V(D^{2} + \lambda I)^{-1}DU^{\top}Y$. Since calculating the inverse of a diagonal matrix is trivial, we see that the complexity of computing $\hat{\beta}_{\lambda}^{R}$ for any λ lies in O(np). Of course, this is after computation of the SVD of X, which lies in $O(np\min(n, p))$.

1.2 v-fold cross validation

Of course, we are still left with the problem of choosing λ in ridge regression. We consider one possible way of doing so, namely v-fold cross validation, which is a general way of selection a good regression method from several competing methods. Here, we assume that our predictors are random, so that we have i.i.d. data pairs (x_i, Y_i) (i = 1, ..., n). Suppose (x^*, Y^*) is a new data pair, independent of

(X,Y) and identically distributed. Ideally, we want to pick λ which minimises the prediction error (averaged over Y^* and x^*)

$$\mathbb{E}\left[\left(Y^* - (x^*)^{\top} \hat{\beta}_{\lambda}^{\mathrm{R}}(X, Y)\right)^2 \mid X, Y\right],$$

where the dependence of $\hat{\beta}_{\lambda}^{R}$ on the training data (X,Y) is made explicit by denoting it $\hat{\beta}_{\lambda}^{R}(X,Y)$.

This is impossible to minimise, but it may be possible to minimise the expected prediction error (averaged over the training data)

$$\mathbb{E}\Big\{\mathbb{E}\Big[\Big(Y^* - (x^*)^\top \hat{\beta}_{\lambda}^{\mathrm{R}}(X, Y)\Big)^2 \mid X, Y\Big]\Big\}. \tag{1}$$

This is still not possible to compute directly, but we estimate it using v-fold cross validation. Split the data into v groups or folds of roughly equal size $(X^{(1)}, Y^{(1)}), \ldots, (X^{(v)}, Y^{(v)})$ and let $(X^{(-k)}, Y^{(-k)})$ denote all data except that in the k-th fold. Then we define

$$CV(\lambda) := \frac{1}{n} \sum_{i=1}^{n} \left[Y_i - x_i^{\mathsf{T}} \hat{\beta}_{\lambda}^{\mathsf{R}} (X^{(-\kappa(i))}, Y^{(-\kappa(i))}) \right]^2,$$

and choose the value of λ that minimises $CV(\lambda)$.

The function $CV(\lambda)$ is called the *out-of-sample error*, since the training data does not include x_i .

Recap 1.8. The tower rule states that for random variables X, Y we have $\mathbb{E}[X] = \mathbb{E}[\mathbb{E}[X \mid Y]]$.

Note that for each i, we have

$$\mathbb{E}\bigg[\Big\{Y_i - x_i^{\intercal} \hat{\beta}_{\lambda}^{\mathrm{R}}(X^{(-\kappa(i))}, Y^{(-\kappa(i))})\Big\}^2\bigg] = \mathbb{E}\bigg[\mathbb{E}\bigg[\Big\{Y_i - x_i^{\intercal} \hat{\beta}_{\lambda}^{\mathrm{R}}(X^{(-\kappa(i))}, Y^{(-\kappa(i))})\Big\}^2 \mid X^{-\kappa(i)}, Y^{-\kappa(i)}\bigg]\bigg].$$

This equals the expected prediction error in eq. (1), except that the training data X, Y are replaced with a smaller data set.

We now have a bias-variance tradeoff in the size of the folds: if v = n (known as "leave-one-out" cross-validation), the estimation will be almost unbiased, but the averaged quantities in $CV(\lambda)$ will be highly correlated which leads to high variance. Typical choices of v are 5 or 10.

Instead of finding the single best λ , we can also aim to find the best weighted combination of λ 's. For example, suppose λ is restricted to a grid $\lambda_1 > \cdots > \lambda_L$. Then we can use any nonnegative least-squares optimization algorithm to minimise

$$\frac{1}{n} \sum_{i=1}^{n} \left[Y_i - \sum_{\ell=1}^{L} w_{\ell} x_i^{\top} \hat{\beta}_{\lambda_{\ell}}^{R} (X^{(-\kappa(i))}, Y^{(-\kappa(i))}) \right]^2,$$

over all $w \in \mathbb{R}^L_{\geq 0}$. This procedure is known as stacking and often outperforms cross-validation.

1.3 The kernel trick

We note that

$$X^{\top}(XX^{\top} + \lambda I) = (X^{\top}X + \lambda I)X^{\top},$$

and multiplying from the left with $(X^{\top}X + \lambda I)^{-1}$ and from the right with $(XX^{\top} + \lambda I)^{-1}$ gives

$$(X^{\top}X + \lambda I)^{-1}X^{\top} = X^{\top}(XX^{\top} + \lambda I)^{-1}.$$

Using this, we see that we can rewrite the fitted values from ridge regression as follows:

$$X\hat{\beta}_{\lambda}^{R} = X(X^{\top}X + \lambda I)^{-1}X^{\top}Y = XX^{\top}(XX^{\top} + \lambda I)^{-1}Y.$$

Two important remarks:

- 1. Computing the LHS of this equation takes roughly $O(np^2 + p^3)$ operations, while computing the RHS takes $O(n^2p + n^3)$ operations (this is because in the LHS we invert an $p \times p$ matrix, while in the RHS we invert a $n \times n$ matrix). Therefore, if $p \gg n$, the RHS can be much cheaper to compute.
- 2. The LHS depends only on the matrix $K = XX^{\top}$ (this matrix is called the *kernel matrix*). Intuitively, since $K_{ij} = \langle x_i, x_j \rangle$, the entries of the kernel matrix show how 'similar' the corresponding predictors are.

Example 1.9. Suppose we have data $(Y_i, z_i)_{i=1,...,n}$ with $z_i = (z_{i1},...,z_{id})^{\top}$, and we believe the following quadratic relation holds:

$$Y_i = \sum_k \sqrt{2} \gamma_k z_{ik} + \sum_{k,\ell} \vartheta_{k\ell} z_{ik} z_{i\ell} + \varepsilon_i.$$

To compute fitted values using ridge regression, we can rewrite this as a linear model $Y = X\beta + \varepsilon$ where

$$\beta = \begin{pmatrix} \gamma_1 \\ \vdots \\ \gamma_d \\ \vartheta_{11} \\ \vartheta_{12} \\ \vdots \\ \vartheta_{dd} \end{pmatrix}, \quad x_i = \begin{pmatrix} \sqrt{2}z_{i1} \\ \vdots \\ \sqrt{2}z_{id} \\ z_{i1}z_{i1} \\ z_{i1}z_{i2} \\ \vdots \\ z_{id}z_{id} \end{pmatrix}.$$

In this case, we have $p = d^2 + d$ variables, which means computing $(X^\top X + \lambda I)^{-1}$ takes $O(d^6)$ operations. In this case, computing $(XX^\top + \lambda I)^{-1}$ is probably easier.

We are still left with the problem of computing $K := XX^{\top}$, which can take $O(n^2p) = O(n^2d^2)$ operations if done naively. However, observe that

$$K_{ij} = x_i^{\top} x_j = 2 \sum_k z_{ik} z_{jk} + \sum_{k,\ell} z_{ik} z_{i\ell} z_{jk} z_{j\ell} = \left(1 + \sum_k z_{ik} z_{jk} \right)^2 - 1 = (1 + z_i^{\top} z_j) - 1.$$

This quantity can be computed in O(d), and therefore K can be computed in $O(n^2d)$ operations: we have a factor d improvement.

The general point of the previous example is that we can bypass the features x_i entirely and instead think directly of $K = XX^{\top}$ where an entry K_{ij} represents similarity between the inputs of the *i*-th and *j*-th samples. This leads to the notion of a kernel in general.

1.4 Kernels

We will assume our inputs x_1, \ldots, x_n live in an abstract space \mathcal{X} .

Definition 1.10. A (positive-definite) kernel is a symmetric map $k: \mathcal{X}^2 \to \mathbb{R}$ such that for all $n \in \mathbb{N}$ and all $x_1, \ldots, x_n \in \mathcal{X}$, the matrix $K \in \mathbb{R}^{n \times n}$ with $K_{ij} = k(x_i, x_j)$ is positive semi-definite.

Proposition 1.11 (Cauchy-Schwarz for kernels). Let k be a kernel and $x, x' \in \mathcal{X}$, then

$$k(x, x')^2 \le k(x, x)k(x', x').$$

Proof. The matrix $\binom{k(x,x)}{k(x',x)} \frac{k(x,x')}{k(x',x')}$ must be positive semi-definite so its determinant must be nonnegative.

In our old models, the data points x_i were vectors in \mathbb{R}^p . Now we try to think of them as points in an abstract space with an associated feature map $\varphi \colon \mathcal{X} \to \mathcal{H}$ (with \mathcal{H} an inner product space), and a kernel k(x, x') gives a measure of similarity between $\varphi(x)$ and $\varphi(x')$. In this case, we have the following:

Proposition 1.12. Let \mathcal{H} be an inner product space, $\varphi \colon \mathcal{X} \to \mathcal{H}$ and define $k(x, x') \coloneqq \langle \varphi(x), \varphi(x') \rangle$. Then k is a kernel.

Proof. We have, for all $x_1, \ldots, x_n \in \mathcal{X}$ and $\alpha \in \mathbb{R}^n$ that

$$\alpha^{\top} K \alpha = \sum_{i,j} K_{ij} \alpha_i \alpha_j = \sum_{i,j} \langle \varphi(x_i), \varphi(x_j) \rangle \alpha_i \alpha_j = \left\| \sum_i \alpha_i \varphi(x_i) \right\|^2 \ge 0.$$

1.4.1 Examples of kernels

The following proposition shows how to make new kernels from old:

Proposition 1.13. Suppose k_1, k_2, \ldots are kernels. Then:

- 1. If $\alpha_1, \alpha_2 \geq 0$ then $\alpha_1 k_1 + \alpha_2 k_2$ is a kernel.
- 2. The pointwise limit of a sequence of kernels is a kernel (if it exists).
- 3. The pointwise product k_1k_2 is a kernel.

Proof. See Example Sheet 1.

Example 1.14. Let us consider some examples of kernels:

- 1. For $\mathcal{X} = \mathbb{R}^p$ we have already seen the linear kernel $k(x, x') = x^{\top} x'$.
- 2. For $\mathcal{X} = \mathbb{R}^p$, the *polynomial kernel* is defined as $k(x, x') = (1 + x^{\top} x')^d$. This is a kernel since it is a power of a sum of two kernels.
- 3. For $\mathcal{X} = \mathbb{R}^p$, the Gaussian kernel is defined by

$$k(x, x') = \exp\left(-\frac{\|x - x'\|_2^2}{2\sigma^2}\right).$$

To show this is a kernel, write k as the pointwise product k_1k_2 where

$$k_1(x, x') = \exp\left(-\frac{\|x\|^2}{2\sigma^2}\right) \exp\left(-\frac{\|x'\|^2}{2\sigma^2}\right), \quad k_2(x, x') = \exp\left(\frac{x^\top x'}{\sigma^2}\right).$$

Clearly k_1 is the kernel induced by the feature map $\varphi(x) = \exp(-\|x\|^2/(2\sigma^2))$, while k_2 can be seen to be a kernel by using the Taylor expansion, which shows that k_2 is a limit of nonnegative linear combinations of kernels.

- 4. For $\mathcal{X} = [0, 1]$, define the Sobolev kernel $k(x, x') = \min(x, x')$. The proof that this is a kernel is on example sheet 1.
- 5. For $\mathcal{X} = \mathcal{P}(\{1,\ldots,p\})$, define the Jaccard kernel

$$k(x, x') = \frac{|x \cap x'|}{|x \cup x'|}$$
 where $0/0 := 1$.

The proof that this is a kernel is on example sheet 1.

1.4.2 Reproducing kernel Hilbert spaces

By proposition 1.12, we see that every feature map $\varphi \colon \mathcal{X} \to \mathcal{H}$ gives rise to a kernel. In the next (important!) theorem, we will see that every kernel is in fact induced by a feature map.

Theorem 1.15. Let k be a kernel, then there exists an inner product space \mathcal{H} and a feature map $\varphi \colon \mathcal{X} \to \mathcal{H}$ such that

$$k(x, x') = \langle \varphi(x), \varphi(x') \rangle$$
 for all $x, x' \in \mathcal{X}$.

Proof. We will construct \mathcal{H} and φ explicitly. First we define the function space

$$\mathcal{H} = \left\{ \sum_{i=1}^{n} \alpha_i k(\cdot, x_i) \mid n \in \mathbb{N}, \alpha_i \in \mathbb{R}, x_i \in \mathcal{X} \right\}.$$

Let $f = \sum_{i=1}^n \alpha_i k(\cdot, x_i)$ and $g = \sum_{j=1}^n \beta_j k(\cdot, x_j')$, then the inner product on \mathcal{H} is given by

$$\langle f, g \rangle = \left\langle \sum_{i=1}^{n} \alpha_i k(\cdot, x_i), \sum_{i=1}^{m} \beta_j k(\cdot, x_j') \right\rangle := \sum_{i=1}^{n} \sum_{j=1}^{m} \alpha_i \beta_j k(x_i, x_j').$$

We define $\varphi \colon \mathcal{X} \to \mathcal{H}$ as $\varphi(x) = k(\cdot, x)$.

We must check that the inner product does not depend on the choice of representation of f and g. For this, note that

$$\sum_{i=1}^{n} \sum_{j=1}^{m} \alpha_i \beta_j k(x_i, x_j') = \sum_{i=1}^{n} \alpha_i g(x_i) = \sum_{j=1}^{m} \beta_j f(x_j'),$$

which holds by symmetry of the kernel. Since $\sum_i \alpha_i g(x_i)$ is independent of the representation of g, while $\sum_j \beta_j f(x'_j)$ is independent of the representation of f, we conclude that the entire expression is independent of both representations.

Secondly, we must verify that the formula $k(x, x') = \langle \varphi(x), \varphi(x') \rangle$ indeed holds. For any $f \in \mathcal{H}, x \in \mathcal{X}$ we have

$$\langle k(\cdot, x), f \rangle = \sum_{i=1}^{n} \alpha_i k(x_i, x) = f(x), \tag{2}$$

i.e., evaluation of a function is a linear functional in \mathcal{H} .

In particular, we have

$$\langle \varphi(x), \varphi(x') \rangle = \langle k(\cdot, x), k(\cdot, x') \rangle = k(x, x').$$

Finally, we must check that $\langle \cdot, \cdot \rangle$ is indeed an inner product. Symmetry and bilinearity are clear. Furthermore, we have

$$\langle f, f \rangle = \sum_{i=1}^{n} \sum_{j=1}^{n} \alpha_i \alpha_j k(x_i, x_j) = \alpha^{\top} K \alpha \ge 0$$

by the fact that k is a kernel. We must now only show that $f \neq 0 \implies \langle f, f \rangle > 0$. For this, note that $\langle \cdot, \cdot \rangle$ is a kernel on \mathcal{H} , so by proposition 1.11 (Cauchy-Schwarz) we have

$$f(x)^2 = \langle k(\cdot, x), f \rangle^2 \le \langle k(\cdot, x), k(\cdot, x) \rangle \langle f, f \rangle,$$

and therefore if f is nonzero anywhere, $\langle f, f \rangle$ must also be nonzero.

While \mathcal{H} constructed in the proof is an inner product space, it is not necessarily a Hilbert space. Let $(f_n) \subseteq \mathcal{H}$ be Cauchy, then by Cauchy-Schwarz for kernels we find

$$f_m(x) - f_n(x) = (f_m - f_n)(x) = \langle k(\cdot, x), f_m - f_n \rangle \le \sqrt{k(x, x)} \|f_n - f_m\|.$$

We can do an analogous computation for $f_n - f_m$ to conclude that $|f_m(x) - f_n(x)| \le \sqrt{k(x,x)} ||f_n - f_m||$, and therefore, if (f_n) is Cauchy, then it converges pointwise to some $f^* : \mathcal{X} \to \mathbb{R}$. We will not prove the following theorem:

Theorem 1.16. The inner product space \mathcal{H} constructed in the proof of theorem 1.15 can be extended to a Hilbert space by adding all pointwise limits f^* of Cauchy sequences $(f_n) \subseteq \mathcal{H}$.

The completion of \mathcal{H} is a special type of Hilbert space:

Definition 1.17. A Hilbert space \mathcal{B} of functions $f: \mathcal{X} \to \mathbb{R}$ is called a reproducing kernel Hilbert space (RKHS) if for all $x \in \mathcal{X}$, there exists $k_x \in \mathcal{B}$ such that

$$f(x) = \langle k_x, f \rangle,$$

i.e., evaluation of functions is a linear functional.

The function $k(x, x') = \langle k_x, k_{x'} \rangle$ is known as the *reproducing kernel* of \mathcal{B} (induced by the feature map $\varphi(x) = k_x$).

If we start with a kernel k, construct the corresponding RKHS \mathcal{B} , then it is easily checked that k is indeed the reproducing kernel of \mathcal{B} .

Example 1.18 (Linear kernel). Let $X = \mathbb{R}^p$ and $k(x, x') = x^{\top} x'$. Then we have

$$\mathcal{H} = \left\{ x \mapsto \sum_{i=1}^{n} \alpha_i x^{\top} x_i = x^{\top} \left(\sum_i \alpha_i x_i \right) \mid \alpha_i \in \mathbb{R}, x_i \in \mathbb{R}^p \right\} = \left\{ x \mapsto x^{\top} \beta \mid \beta \in \mathbb{R}^p \right\},$$

and if $f(x) = x^{\top}\beta$, $g(x) = x^{\top}\beta'$, then

$$\langle f, g \rangle = k(\beta, \beta') = \beta^{\top} \beta' \text{ so } ||f||_{\mathcal{H}} = ||\beta||_{2}.$$

1.4.3 The representer theorem

We started with data (Y_i, x_i) in a model where our fitted values depend only on $XX^{\top} = K$, where K is the kernel matrix derived from the kernel $k(x_i, x_j) = \langle x_i, x_j \rangle = x_i^{\top} x_j$. Suppose we change the kernel to some \tilde{k} , and replace K by \tilde{K} accordingly. Then by theorem 1.15 this is equivalent by replacing our predictors x_i by $\varphi(x_i)$ for some feature map φ .

Recall that the ridge regression objective function is given by

$$\min_{\beta \in \mathbb{R}^p} \sum_{i=1}^n (Y_i - x_i^\top \beta)^2 + \lambda \|\beta\|^2.$$

Let \mathcal{H} be the RKHS of the linear kernel, then the above is equivalent to

$$\min_{f \in \mathcal{H}} \sum_{i=1}^{n} (Y_i - f(x_i))^2 + \lambda \|f\|^2.$$
 (3)

We will show that in general, if we change our predictors x_i to $\varphi(x_i)$, and update the corresponding RKHS in eq. (3), then eq. (3) solves our new problem. Note that there are no references to the feature map φ in eq. (3).

Theorem 1.19 (Representer theorem). *Let:*

- $(Y_i, x_i)_{i=1}^n$ be our data with $x_i \in \mathcal{X}$;
- $c: \mathbb{R}^n \times \mathcal{X}^n \times \mathbb{R}^n \to \mathbb{R}$ be an arbitrary loss function;
- $J: \mathbb{R}_{\geq 0} \to \mathbb{R}$ strictly increasing;
- k a kernel with RKHS \mathcal{H} and kernel matrix $K_{ij} = k(x_i, x_j)$.

Then $\hat{f} \in \mathcal{H}$ minimises

$$Q_1(f) := c(Y, x_1, \dots, x_n, f(x_1), \dots, f(x_n))$$

if and only if $\hat{f}(\cdot) = \sum_{i=1}^{n} \hat{\alpha}_i k(\cdot, x_i)$, where $\hat{\alpha} \in \mathbb{R}^n$ minimises

$$Q_2(\alpha) := c(Y, x_1, \dots, x_n, K\alpha) + J(\alpha^\top K\alpha).$$

Proof. Let $f \in \mathcal{H}$ and $U = \operatorname{Span}\{k(\cdot, x_1), \dots, k(\cdot, x_n)\} \subseteq \mathcal{H}$, then U is closed since it is finite-dimensional so we can write f = u + v with $u \in U, v \in U^{\perp}$. Note that

$$f(x_i) = \langle k(\cdot, x_i), u + v \rangle = \langle k(\cdot, x_i), u \rangle = u(x_i) \quad (i = 1, \dots, n),$$

so minimisation of Q_1 is equivalent to minimisation of

$$c(Y, x_1, \dots, x_n, u(x_1), \dots, u(x_n)) + J(\|u\|^2 + \|v\|^2)$$
 (4)

w.r.t. $u \in U$, $v \in U^{\perp}$. Now, since J is strictly increasing, any minimiser (u, v) of eq. (4) will satisfy v = 0.

Write $u(\cdot) = \sum_{i=1}^{n} \alpha_i k(\cdot, x_i) \in U$, then clearly

$$\begin{bmatrix} u(x_1) \\ \vdots \\ u(x_n) \end{bmatrix} = \begin{bmatrix} \sum_i \alpha_i k(x_1, x_i) \\ \vdots \\ \sum_i \alpha_i k(x_n, x_i) \end{bmatrix} = K\alpha \quad \text{and} \quad \|u\|^2 = \sum_{i,j} \alpha_i \alpha_j k(x_i, x_j) = \alpha^\top K\alpha.$$

This shows that minimisation of eq. (4) is equivalent to minimisation of

$$c(Y, x_1, \dots, x_n, K\alpha) + J(\alpha^\top K\alpha) = Q_2(\alpha),$$

which completes the proof.

Example 1.20. In ridge regression, the representer theorem tells us that minimising eq. (3) is equivalent to minimising $\|Y - K\alpha\|^2 + \lambda \alpha^\top K\alpha$, and indeed it can be shown (example sheet 1) that the minimiser satisfies $K\hat{\alpha} = K(K + \lambda I)^{-1}Y = X\hat{\beta}^R_{\lambda}$. Since \mathcal{H} may be infinite-dimensional, this gives a way to rewrite an infinite-dimensional optimization problem to a finite-dimensional optimisation problem.

1.5 Kernel ridge regression

We have now defined kernel ridge regression and shown how the estimator may be computed, but we have yet to assess its predictive performance. We consider the model

$$Y_i = f^0(x_i) + \varepsilon_i, \quad \mathbb{E}[\varepsilon] = 0, \text{Var}[\varepsilon] = \sigma^2 I,$$

where we assume $f^0 \in \mathcal{H}$ where \mathcal{H} is an RKHS with reproducing kernel k. By scaling the equation on both sides we may assume $||f^0|| \leq 1$ (note that this changes $\operatorname{Var}[\varepsilon]$). Let K be the kernel matrix with eigenvalues $d_1 \geq \cdots \geq d_n \geq 0$, and let $\hat{f} = \hat{f}_{\lambda}$ be the estimated regression function, so

$$\hat{f}_{\lambda} = \underset{f \in \mathcal{H}}{\operatorname{arg \, min}} \left(\sum_{i=1}^{n} (Y_i - f(x_i))^2 + \lambda ||f||^2 \right).$$

Theorem 1.21. The mean squared prediction error (MSPE) satisfies

$$\frac{1}{n}\mathbb{E}\left[\sum_{i=1}^{n} \left(f^{0}(x_{i}) - \hat{f}_{\lambda}(x_{i})\right)^{2}\right] \leq \frac{\sigma^{2}}{n} \sum_{i=1}^{n} \frac{d_{i}^{2}}{(d_{i} + \lambda)^{2}} + \frac{\lambda}{4n}$$

$$\leq \frac{\sigma^{2}}{n\lambda} \sum_{i=1}^{n} \min\left(\frac{d_{i}}{4}, \lambda\right) + \frac{\lambda}{4n}.$$

Proof. The representer theorem tells us that $(\hat{f}(x_1), \dots, \hat{f}(x_n))^{\top} = K(K + \lambda I)^{-1}Y$. By projecting f^0 onto Span $k(\cdot, x_1), \dots, k(\cdot, x_n)$ it is easily seen that there exists $\alpha \in \mathbb{R}^n$ such that

$$(f^0(x_1), \dots, f^0(x_n))^\top = K\alpha \text{ and } \|f^0\|^2 \ge \alpha^\top K\alpha.$$

Write $K = UDU^{\top}$ where $D_{ii} = d_i$ and define $\vartheta := U^{\top}K\alpha = DU^{\top}\alpha$ (note $U\vartheta = K\alpha$ and note that $d_i = 0 \implies \vartheta_i = 0$). Then we have

$$\mathbb{E}\left[\sum_{i=1}^{n} \left(f^{0}(x_{i}) - \hat{f}_{\lambda}(x_{i})\right)^{2}\right] = \mathbb{E}\left\|K\alpha - K(K + \lambda I)^{-1}Y\right\|^{2}$$

$$= \mathbb{E}\left\|K(K + \lambda I)^{-1}(U\vartheta + \varepsilon) - U\vartheta\right\|^{2}$$

$$= \mathbb{E}\left\|DU^{\top}(UDU^{\top} + \lambda I)^{-1}(U\vartheta + \varepsilon) - \vartheta\right\|^{2}$$

$$= \mathbb{E}\left\|D(D + \lambda I)^{-1}(\vartheta + U^{\top}\varepsilon) - \vartheta\right\|^{2}$$

$$= \mathbb{E}\left\|\left[D(D + \lambda I)^{-1} - I\right]\vartheta + D(D + \lambda I)^{-1}U^{\top}\varepsilon\right\|^{2}$$

$$= \left\|\left[D(D + \lambda I)^{-1} - I\right]\vartheta\right\|^{2} + \mathbb{E}\left\|D(D + \lambda I)^{-1}U^{\top}\varepsilon\right\|^{2}$$

Note that the cross-term in the final equality disappears since it is the expectation of a linear combination of ε , which is 0.

The first term is a deterministic quantity, which equals

$$\left\| \left[D(D+\lambda I)^{-1} - I \right] \vartheta \right\|^2 = \sum_{i=1}^n \left(\left(\frac{d_i}{d_i + \lambda} - 1 \right) \vartheta \right)^2 = \sum_{i=1}^n \frac{\lambda^2 \vartheta_i^2}{(d_i + \lambda)^2}.$$

Let D^+ be the diagonal matrix with $D_{ii}^+ = d_i^{-1}$ if $d_i \neq 0$ and 0 else. Then we have

$$\sum_{i:d_i>0} \frac{\vartheta_i^2}{d_i} = \vartheta^\top D^+ \vartheta = \alpha^\top K U D^+ U^\top K \alpha = \alpha^\top U D D^+ D U^\top \alpha = \alpha^\top U D U^\top \alpha = \alpha^\top K \alpha \ge 1,$$

Using this, we can bound the first term by

$$\sum_{i=1}^n \frac{\lambda^2 \vartheta_i^2}{(d_i + \lambda)^2} = \sum_{i: d_i \neq 0} \frac{\vartheta_i^2}{d_i} \frac{d_i \lambda^2}{(d_i + \lambda)^2} \le \max_i \frac{d_i \lambda^2}{(d_i + \lambda)^2} \le \max_i \frac{d_i \lambda^2}{4\lambda d_i} \le \frac{\lambda}{4}.$$

To compute the second term, we use the *trace trick*:

$$\mathbb{E} \|D(D+\lambda I)^{-1}U^{\top}\varepsilon\|^{2} = \mathbb{E} \Big[(D(D+\lambda I)^{-1}U^{\top}\varepsilon)^{\top} (D(D+\lambda I)^{-1}U^{\top}\varepsilon) \Big]$$

$$= \mathbb{E} \Big[\operatorname{tr} \big[D(D+\lambda I)^{-1}U^{\top}\varepsilon\varepsilon^{\top}U(D+\lambda I)^{-1}D \big] \Big]$$

$$= \operatorname{tr} \big[D(D+\lambda I)^{-1}U^{\top}\mathbb{E}(\varepsilon\varepsilon^{\top})U(D+\lambda I)^{-1}D \big]$$

$$= \sigma^{2} \operatorname{tr} \big[D^{2}(D+\lambda I)^{-2} \big] = \sigma^{2} \sum_{i=1}^{n} \frac{d_{i}^{2}}{(d_{i}+\lambda)^{2}}.$$

This gives our first inequality MSPE $\leq \frac{\sigma^2}{n} \sum_{i=1}^n \frac{d_i^2}{(d_i+\lambda)^2} + \frac{\lambda}{4n}$. For the second inequality, note that $\frac{d_i^2}{(d_i+\lambda)^2} \leq 1$ and also $\frac{d_i^2}{(d_i+\lambda)^2} \leq \frac{d_i^2}{4d_i\lambda} = \frac{d_i}{4\lambda}$. Therefore, we can write $\frac{d_i^2}{(d_i+\lambda)^2} \leq \frac{1}{\lambda} \min(\lambda, \frac{d_i}{4})$ which proves the second inequality.

We can now ask ourselves if kernel ridge regression is optimal in any sense. To this end, we define $\hat{\mu}_i := d_i/n$ and $\lambda_n := \lambda/n$. Then we can rewrite the upper bound from our precious theorem as

$$MSPE \leq \frac{\sigma^2}{n\lambda_n} \sum_{i=1}^n \min(\frac{\hat{\mu}_i}{4}, \lambda_n) + \frac{\lambda_n}{4} =: \delta_n(\lambda_n).$$

Now, instead of taking $x_1, \ldots, x_n \in \mathcal{X}$ to be fixed, assume that they are i.i.d. random variables (then K is random, so the $\hat{\mu}_i$ are random as well), and taking an expectation on both sides yields

$$\mathbb{E}[MSPE] \leq \mathbb{E}[\delta_n(\lambda_n)].$$

We want to bound $\delta_n(\lambda_n)$ by some function of n. For this we will use Mercer's theorem:

Theorem 1.22 (Mercer). Let $\mathcal{X} = [a,b]$. If k is a continuous kernel on \mathcal{X} , then there exists an orthonormal basis (e_i) of $L^2[a,b]$ such that

$$k(x, x') = \sum_{j=1}^{\infty} \mu_j e_j(x) e_j(x')$$

as well as a sequence of nonnegative eigenvalues (μ_i) such that

$$\int_a^b k(x, x')e_j(x) dx = \mu_j e_j(x).$$

Applying Mercer's theorem to k, it can be proved that for some C > 0 we have

$$\mathbb{E}\left[\frac{1}{n}\sum_{i=1}^{n}\min(\frac{\hat{\mu}_{i}}{4},\lambda_{n})\right] \leq \frac{C}{n}\sum_{i=1}^{\infty}\min(\frac{\mu_{i}}{4},\lambda_{n}).$$

Example 1.23. Let k be the Sobolev kernel. Then the eigenvalues satisfy

$$\frac{\mu_j}{4} = \frac{1}{\pi^2 (2j-1)^2},$$

and with some calculations it can be shown that $\sum_{i=1}^{\infty} \min(\frac{\mu_i}{4}, \lambda_n) = O(\sqrt{\lambda_n})$ as $\lambda_n \to 0$. Combining this with the rest of the bound shows

$$\mathbb{E}[\delta_n(\lambda_n)] = O\left(\frac{\sigma^2}{n}\lambda_n^{-1/2} + \lambda_n\right).$$

The optimal scaling $\lambda_n \sim (\sigma^2/n)^{2/3}$ gives an error rate of order $(\sigma^2/n)^{2/3}$.

It has been shown that this is in fact the optimal rate for estimating a function $f^0 \in \mathcal{H}$ (up to multiplicative constants).

1.6 Other kernel machines

1.6.1 The support vector machine

Consider a classification problem with data $x_1, \ldots, x_n \in \mathbb{R}^p$ and binary response $Y_i \in \{\pm 1\}$. Furthermore assume that x_1, \ldots, x_n are separated by a hyperplane through the origin, that is, for some $\beta \in \mathbb{R}^p$ we have $Y_i x_i^{\top} \beta > 0$ for all i.

To choose between different planes that separate the classes, we pick the hyperplane with the highest margin: that is, we compute

$$\max_{\beta \in \mathbb{R}^p, M \ge 0} M \quad \text{such that } \frac{Y_i x_i^\top \beta}{\|\beta\|} \ge M \quad (i = 1, \dots, n).$$

When the classes are not separable, we can penalise according to the distance of a point "over the margin". The penalty should be 0 if x is on the correct side and should equal the distance over the boundary otherwise, measured in units of M. In this case, by considering minimisation of $1/M^2$ instead of maximisation of M we find

$$\underset{M \geq 0, \beta \in \mathbb{R}^p}{\operatorname{arg\,min}} \frac{1}{M^2} + \lambda \sum_{i=1}^n \left(1 - \frac{Y_i x_i^\top \beta}{\|\beta\| M} \right)_+.$$

Since the above equation is independent of the norm of β , we may set $\|\beta\| = 1/M$ and remove M from the equation entirely, thus obtaining

$$\underset{\beta \in \mathbb{R}^p}{\operatorname{arg\,min}} \|\beta\|^2 + \lambda \sum_{i=1}^n (1 - Y_i x_i^{\top} \beta)_+.$$

Replacing λ by $\frac{1}{\lambda}$ and multiplying the equation with λ we obtain

$$\underset{\beta \in \mathbb{R}^p}{\operatorname{arg\,min}} \sum_{i=1}^n (1 - Y_i x_i^{\top} \beta)_+ + \lambda \|\beta\|^2.$$

Finally, if we lift the restriction that the hyperplane must go through the origin, we find the objective function

$$(\hat{\mu}, \hat{\beta}) = \underset{\mu \in \mathbb{R}, \beta \in \mathbb{R}^p}{\min} \sum_{i=1}^{n} (1 - Y_i(x_i^{\top} \beta + \mu))_+ + \lambda \|\beta\|^2.$$
 (5)

The solution to this problem is known as the *support vector classifier*.

Introducing kernels Now we introduce kernels to the problem. Letting k be the linear kernel and \mathcal{H} the corresponding RKHS, we can rewrite eq. (5) as

$$(\hat{\mu}, \hat{f}) = \underset{\mu \in \mathbb{R}, f \in \mathcal{H}}{\arg \min} \sum_{i=1}^{n} (1 - Y_i^{\top} (f(x_i) + \mu))_+ + \lambda ||f||^2.$$

Suppose we change \mathcal{H} to a different RKHS with a different kernel k. By a variant of the representer theorem (Example Sheet 1, question 9) the solution is equivalent to

$$(\hat{\mu}, \hat{\alpha}) = \underset{\mu \in \mathbb{R}, \alpha \in \mathbb{R}^n}{\operatorname{arg \, min}} \sum_{i=1}^n (1 - Y_i(K_i^{\top} \alpha + \mu))_+ + \lambda \alpha^{\top} K \alpha,$$

and the solution to the above problem is called the *support vector machine*. Predictions at a new point x^* are given by

$$\operatorname{sign}(\hat{\mu} + \sum_{i=1}^{n} \hat{\alpha}_i k(x^*, x_i)) \in \pm 1.$$

Remark. Note that the support vector machine can correspond to a nonlinear boundary: the boundary is a hyperplane in the corresponding RKHS, but can be nonlinear in the space \mathcal{X} where the data points live. This also makes it easy to overfit.

1.6.2 Logistic regression

Recap 1.24. The standard logistic regression model is motivated by assuming

$$\log\left(\frac{\mathbb{P}(Y_i=1)}{\mathbb{P}(Y_i=-1)}\right) = x_i^{\top} \beta^0.$$

The maximum-likelihood estimator $\hat{\beta}$ is given by (Example sheet 2, question 4)

$$\hat{\beta} := \operatorname*{arg\,min}_{\beta \in \mathbb{R}^p} \sum_{i=1}^n \log \left(1 + \exp \left(-Y_i x_i^\top \beta \right) \right).$$

Just like in ridge regression, we will try to lower the variance of $\hat{\beta}$ by penalising large values, i.e.,

$$\hat{\beta}_{\lambda} := \underset{\beta \in \mathbb{R}^p}{\operatorname{arg\,min}} \sum_{i=1}^n \log(1 + \exp(-Y_i x_i^{\top} \beta)) + \lambda \|\beta\|^2$$
$$= \underset{f \in \mathcal{H}}{\operatorname{arg\,min}} \sum_{i=1}^n \log(1 + \exp(-Y_i f(x_i))) + \lambda \|f\|^2,$$

where \mathcal{H} is the RKHS corresponding to the linear kernel.

Note that this is exactly the same as the support vector machine, but with a different loss function. The loss function in the support vector machine is called "hinge loss", while the loss function in logistic regression is called "logistic loss". The main differences are:

- 1. The logistic loss function increases much quicker for incorrectly classified data.
- 2. Even when a point is correctly classified, the logistic loss function decreases the further the datapoint is from the boundary. The hinge loss function will always be 0 if a point is correctly classified and lies outside of the margin.

The second point also explains the name *support vector machine*: the only points that contribute to the loss function are those that either violate the boundary of are within the margin (the *support vectors*), and perturbing any of the other points does not change the boundary. Because of this property, SVM's tend to be more stable.

1.7 Large scale kernel machines

We have seen that kernels give very flexible regression and classification estimators, and that the parameter λ helps us with the bias-variance tradeoff. Also, when $p \gg n$, using the kernel matrix saves computational effort.

We now turn to the case $n \gg p$. In this case, both kernel ridge regression and the SVM are very expensive to compute: the former has cost $O(n^3)$, while the latter can be computed using an iterative algorithm where every iteration has cost $O(n^2)$.

One approach to speed this process up is that of random feature maps: we develop a random map $\hat{\varphi} \colon \mathcal{X} \to \mathbb{R}^b$ (with b small) such that $\mathbb{E}[\hat{\varphi}(x)^\top \hat{\varphi}(x')] = k(x, x')$ for all $x, x' \in \mathcal{X}$. Then we consider L i.i.d. copies of $\hat{\varphi}$ (denotes $\hat{\varphi}_1, \ldots, \hat{\varphi}_n$), and we consider the feature map

$$\psi(x) = L^{-1/2}(\hat{\varphi}_1(x), \dots, \hat{\varphi}_L(x))^{\top} \in \mathbb{R}^{Lb}$$
, so that $\psi(x)^{\top} \psi(x') = \frac{1}{L} \sum_{i=1}^{L} \hat{\varphi}_i(x)^{\top} \hat{\varphi}_i(x')$.

Letting Φ be the matrix with rows $\psi(x_1)^{\top}, \ldots, \psi(x_n)^{\top}$, we find that $\mathbb{E}(\Phi\Phi^{\top})_{ij} = k(x_i, k_j)$, and that $\operatorname{Var}(\Phi\Phi^{\top})_{ij}$ decreases with speed L^{-1} . So $\Phi\Phi^{\top}$ is a good approximation for K when L is moderately large.

One specific example for shift-invariant kernels (k(x, x') = h(x - x')) for all $x, x' \in \mathcal{X}$ is based on the work of Rahimi and Recht.

Theorem 1.25 (Bochner). Let $k: \mathbb{R}^p \times \mathbb{R}^p \to \mathbb{R}$ be a continuous kernel. Then k is shift-invariant if and only if, for some c > 0 and some distribution F on \mathbb{R}^p , it holds for all $W \sim F$ that

$$k(x, x') = c\mathbb{E}\left[\exp\left(i(x - x')^{\top}W\right)\right] = c\mathbb{E}\left[\cos\left((x - x')^{\top}W\right)\right].$$

To make use of this theorem, we let $u \sim U[-\pi, \pi]$, and it can then be computed (using $u \stackrel{d}{=} -u$ and $\cos(u) \stackrel{d}{=} \sin(u)$ that

$$2\mathbb{E}[\cos(x+u)\cos(y+u)] = \cos(x-y).$$

Therefore, given a shift-invariant kernel k with associated distribution F, we define for $W \sim F$ the feature map

$$\hat{\varphi}(x) = \sqrt{2c}\cos(W^{\top}x + u) \in \mathbb{R},$$

so that

$$\mathbb{E}\hat{\varphi}(x)\hat{\varphi}(x') = 2c\mathbb{E}\left[\mathbb{E}[\cos(W^{\top}x + u)\cos(W^{\top}x' + u) \mid W]\right] = c\mathbb{E}\cos((x - x^{\top})W) = k(x, x').$$

Example 1.26. Let k be the Gaussian kernel $k(x, x') = e^{-\|x - x'\|/(2\sigma^2)}$. Let $W \sim N(0, \sigma^{-2}I)$, then using the characteristic function of a normal distribution we find that

$$\mathbb{E}\Big(e^{i(x-x')^{\top}W}\Big) = \exp\Big(-\|x-x'\|^2/(2\sigma^2)\Big) = k(x,x'),$$

so we can set c=1, i.e., define the feature map $\hat{\varphi}(x)=\sqrt{2}\cos(W^{\top}x+u)$.

2 The Lasso

2.1 Model selection

We go back to the linear model $Y = X\beta^0 + \varepsilon$ with $\mathbb{E}[\varepsilon] = 0$ and $Var(\varepsilon) = \sigma^2 I$. Using the trace trick, one can easily compute that the MSPE of the OLS estimator is given by

$$\frac{1}{n} \mathbb{E} \left\| X \beta^0 - X \hat{\beta}^{\text{OLS}} \right\|_2^2 = \frac{\sigma^2 p}{n}.$$

Defining $S = \{k \mid (\beta^0)_k \neq 0\}$, there is often reason to assume that S is small, i.e., $s := |S| \ll p$. If we could fit a model using only the variables in S, the MSPE would be much $\frac{\sigma^2 s}{n} \ll \frac{\sigma^2 p}{n}$.

Best subset selection A natural way to find S is to consider all possible subsets of $\{1, \ldots, p\}$, and pick the best regression procedure using, for example, cross-validation. However, this can become computationally infeasible for moderately large p (say $p \approx 10$).

Forward selection This is a greedy way of performing best subset regression. Given a target model size m, we first compute the intercept-only model M_0 , and then one-by-one add the predictor variable that reduces the residual sum of squares the most, until we have a model with m variables.

2.2 Lasso estimator

The Least absolute shrinkage and selection operator or Lasso is given by

$$(\hat{\mu}_{\lambda}^{\mathrm{L}}, \hat{\beta}_{\lambda}^{\mathrm{L}}) \coloneqq \underset{(\mu, \beta) \in \mathbb{R} \times \mathbb{R}^p}{\arg \min} \frac{1}{2n} \|Y - \mu \mathbf{1} - X\beta\|_2^2 + \lambda \|\beta\|_1.$$

As with ridge regression, we usually centre and scale the matrix X, as well as centre the responses Y, in which case we find

$$\hat{\beta}_{\lambda}^{\mathrm{L}} = \operatorname*{arg\,min}_{\beta \in \mathbb{R}^p} \frac{1}{2n} \|Y - X\beta\|_2^2 + \lambda \|\beta\|_1.$$

The main difference between the lasso estimator and the ridge regression estimator is that it is likely that the lasso estimator has some zero components. This means that the lasso estimator also estimates which variables are relevant.

We have the following:

Theorem 2.1 (Slow rate). Assume X has centred and scaled columns, and assume that Y has been centred, so $Y = X\beta^0 + \varepsilon - \bar{\varepsilon} \mathbf{1}$. Let A > 0 and suppose

$$\lambda = A\sigma\sqrt{\frac{\log(p)}{n}}.$$

Let $\hat{\beta} = \hat{\beta}_{\lambda}^{L}$, then with probability at least $1 - 2p^{-(A^{2}/2-1)}$ we have that the MSPE satisfies

$$\frac{1}{n} \left\| X(\beta^0 - \hat{\beta}) \right\|_2^2 \leq 4\lambda \big\| \beta^0 \big\|_1 = 4A\sigma \sqrt{\frac{\log(p)}{n}} \big\| \beta^0 \big\|_1.$$

Proof. By definition we have

$$\frac{1}{2n} \left\| Y - X \hat{\beta} \right\|_{2}^{2} + \lambda \left\| \hat{\beta} \right\|_{1} \leq \frac{1}{2n} \left\| Y - X \beta^{0} \right\|_{2}^{2} + \lambda \left\| \beta^{0} \right\|_{1},$$

and rearranging the terms gives

$$\frac{1}{2n} \left\| X(\beta^0 - \hat{\beta}) \right\|_2^2 \le \frac{1}{n} \varepsilon^\top X(\hat{\beta} - \beta^0) + \lambda \left\| \beta^0 \right\|_1 - \lambda \left\| \hat{\beta} \right\|_1.$$

By Hölder's inequality we have $\left|\varepsilon^{\top}X(\beta^{0}-\hat{\beta})\right| \leq \left\|X^{\top}\varepsilon\right\|_{\infty} \left\|\hat{\beta}-\beta^{0}\right\|_{1}$. Define the event $\Omega = \left\{\left\|X^{\top}\varepsilon\right\|_{\infty}/n \leq \lambda\right\}$, then conditional on Ω we find

$$\frac{1}{n} \| X(\beta^0 - \hat{\beta}) \|_2^2 \le 2\lambda \Big(\| \beta^0 - \hat{\beta} \|_1 + \| \beta^0 \|_1 - \| \hat{\beta} \|_1 \Big) \le 4\lambda \| \beta^0 \|_1.$$

In lemma 2.5, we will show that $\mathbb{P}(\Omega) \geq 1 - 2p^{-(A^2/2-1)}$, which completes the proof.

2.3 Concentration inequalities

Let W be any random variable and $\varphi \colon \mathbb{R} \to [0, \infty)$ strictly increasing. Then by Markov's inequality we have

$$\mathbb{P}(W \ge t) = \mathbb{P}(\varphi(W) \ge \varphi(t)) \le \frac{\mathbb{E}[\varphi(W)]}{\varphi(t)}.$$

Plugging in $\varphi(x) = e^{\alpha x}$ (for some $\alpha > 0$), we get

$$\mathbb{P}(W \ge t) \le e^{-\alpha t} \mathbb{E}[e^{\alpha W}] = e^{-\alpha t} M_W(\alpha).$$

Now we can take the infimum over all α on the right-hand side, and we get what is called the *Chernoff* bound:

$$\mathbb{P}(W \ge t) \le \inf_{\alpha > 0} e^{-\alpha t} M_W(\alpha).$$

Definition 2.2. A random variable W with mean μ is called *sub-Gaussian with parameter* $\sigma > 0$ or σ -sub-Gaussian if

$$M_{W-\mu} \leq M_{N(0,\sigma^2)} \quad \text{or equivalently} \quad \mathbb{E}[e^{\alpha(W-\mu)}] \leq e^{\alpha^2\sigma^2/2} \text{ for all } \alpha \in \mathbb{R}.$$

We need the following lemma, which characterises an important class of sub-Gaussian random variables:

Lemma 2.3 (Hoeffding). If W is a mean-zero random variable which takes values in [a,b], then W is sub-Gaussian with parameter (b-a)/2.

By Chernoff bounding, we obtain for a σ -sub-Gaussian random variable W that

$$\mathbb{P}(W - \mu \ge t) \le e^{-t^2/(2\sigma^2)}.$$

Proposition 2.4. Let W_1, \ldots, W_n be independent, mean-zero random variables where W_i is σ_i -sub-Gaussian. For any $\gamma \in \mathbb{R}^n$, the random variable $\gamma^\top W$ is sub-Gaussian with parameter $\left(\sum_i \sigma_i^2 \gamma_i^2\right)^{1/2}$.

Proof. Since the W_i are independent we have for all $\alpha \in \mathbb{R}$ that

$$\mathbb{E}[e^{\alpha \sum_i \gamma_i W_i}] = \prod_{i=1}^n \mathbb{E}[e^{\alpha \gamma_i W_i}] \le \prod_{i=1}^n e^{\alpha^2 \gamma_i^2 \sigma_i^2/2} = e^{\alpha^2 \sum_i \gamma_i^2 \sigma_i^2/2}.$$

Lemma 2.5. Suppose $\varepsilon_1, \ldots, \varepsilon_n$ are independent mean-zero σ -sub-Gaussian random variables and let $\lambda := A\sigma\sqrt{\log(p)/n}$. Then

$$\mathbb{P}\left(\frac{\left\|X^{\top}\varepsilon\right\|_{\infty}}{n} \le \lambda\right) \ge 1 - 2p^{-(A^2/2 - 1)}.$$

Proof. We have

$$\mathbb{P}\bigg(\frac{\left\|X^{\top}\varepsilon\right\|_{\infty}}{n}>\lambda\bigg)=\mathbb{P}\bigg(\bigcup_{i}\frac{\left|X_{i}^{\top}\varepsilon\right|}{n}>\lambda\bigg)\leq\sum_{i}\mathbb{P}\bigg(\frac{\left|X_{i}^{\top}\varepsilon\right|}{n}>\lambda\bigg).$$

By the previous proposition, both $X_i^{\mathsf{T}}\varepsilon$ and $-X_i^{\mathsf{T}}\varepsilon$ are sub-Gaussian with parameter σ/\sqrt{n} , so we have

$$\sum_i \mathbb{P}\Bigg(\frac{\left|X_i^\top \varepsilon\right|}{n} > \lambda\Bigg) \leq 2p \exp\bigg(\frac{-\lambda^2}{2(\sigma/\sqrt{n})^2}\bigg) = 2p \exp\Big(-A^2 \log(p)/2\Big) = 2p^{-(A^2/2 - 1)}.$$

We now move to finding tail bounds for products of sub-Gaussian random variables. To this end, we consider Bernstein's in equality.

Definition 2.6. A random variable W with mean μ is said to satisfy Bernstein's condition with parameters $\sigma > 0, b > 0$ if

$$\mathbb{E}(|W - \mu|^k) \le \frac{1}{2}k!\sigma^2b^{k-2} \qquad k = 2, 3, \dots$$

Proposition 2.7 (Bernstein's inequality). Let W_1, W_2, \ldots be independent random variables with mean μ , where each W_i satisfies Bernstein's condition with parameter (σ, b) . Then

$$1. \ \mathbb{E}[\exp(\alpha(W_i - \mu))] \leq \exp\left(\frac{\alpha^2 \sigma^2}{2(1 - b|\alpha|)}\right) \ for \ all \ |\alpha| < b^{-1};$$

2.
$$\mathbb{P}(\overline{W} - \mu \ge t) \le \exp\left(\frac{-nt^2}{2(\sigma^2 + bt)}\right)$$
 for all $t \ge 0$.

Proof. See notes from topics in statistical theory.

Lemma 2.8. Let W, Z be mean-zero, sub-Gaussian random variables with parameters σ_W, σ_Z . Then WZ satisfies the Bernstein condition with parameters $(8\sigma_W\sigma_Z, 4\sigma_W\sigma_Z)$.

Proof. Using the tail probability formula for expectations one can prove that any mean-zero sub-Gaussian random variable X satisfies $\mathbb{E}[X^{2k}] \leq 2^{k+1} \sigma_X^{2k} k!$.

Furthermore for any random variable Y we have

$$\begin{split} \mathbb{E}|Y - \mathbb{E}Y|^k &\leq \mathbb{E}(|Y| + |\mathbb{E}Y|)^k \\ &= \mathbb{E}\sum_{t=0}^k \binom{k}{t} |Y|^t |\mathbb{E}Y|^{k-t} \\ &\stackrel{*}{\leq} \sum_{t=0}^k \binom{k}{t} \mathbb{E}(|Y|^t) \mathbb{E}(|Y|^{k-t}) \\ &\stackrel{*}{\leq} \sum_{t=0}^k \binom{k}{t} \mathbb{E}(|Y|^k)^{t/k} \mathbb{E}(|Y|^k)^{k-t/k} \\ &= \mathbb{E}(|Y|^k) \sum_{t=0}^k \binom{k}{t} = 2^k \mathbb{E}|Y|^k, \end{split}$$

where both equations * are applications of Jensen's inequality.

From the above bound, we obtian

$$\mathbb{E}|WZ - \mathbb{E}WZ|^{k} \le 2^{k} \mathbb{E}(|W|^{k}|Z|^{k}) \stackrel{\text{CS}}{\le} 2^{k} (\mathbb{E}W^{2k})^{1/2} (\mathbb{E}Z^{2k})^{1/2} \le 2^{2k+1} \sigma_{W}^{k} \sigma_{Z}^{k} k!$$
$$= \frac{k!}{2} (8\sigma_{W}\sigma_{Z})^{2} (4\sigma_{W}\sigma_{Z})^{k-2},$$

which proves the claim.

2.3.1 Optimisation theory and convex analysis

Proposition 2.9. Let $f: \mathbb{R}^d \to \mathbb{R}$ be twice continuously differentiable. Then f is convex (resp. strictly convex) if and only if its Hessian is positive semi-definite (resp. positive definite) for all $x \in \mathbb{R}^d$.