MATH 590: Meshfree Methods Machine Learning

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Outline



- 2 Radial Basis Function Networks
- Classification with Support Vector Machines Theory
- Olassification with Support Vector Machines Practice
- 5 Support Vector Regression



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Given data

$$\{(\boldsymbol{x}_i, y_i)\}_{i=1}^N, \qquad \boldsymbol{x}_i \in \Omega \subset \mathbb{R}^d, \quad y_i \in \mathbb{R},$$

find a function s that predicts, for a previously unobserved \boldsymbol{x} value,

 $s(\mathbf{x}) \approx \mathbf{y}.$



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Scattered data interpolation: Construct *s* as a linear combination of "shifts" of kernels such that $||\boldsymbol{s} - \boldsymbol{y}|| = 0$, where $\boldsymbol{y} = (y_1, \dots, y_N)^T$ and $\boldsymbol{s} = (s(\boldsymbol{x}_1), \dots, s(\boldsymbol{x}_N))^T$.



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Truncated Mercer series: Construct *s* as a linear combination of only M < N eigenfunctions such that min_{*s*} ||*s* - *y*||.

We now consider ill-posed problems

A problem may be ill-posed if, e.g.,

- we don't have enough data to capture the complexity of the model,
- we don't have enough complexity in our model to match the data,
- we don't want to match all the complexity of the data because some of it might be due to measurement errors.



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In such cases our earlier approaches need to be modified and one typically solves the data fitting problem via a regularization approach. We now

- give a overview to such a general regularization strategy to fitting data,
- look at RBF network regression,
- support vector machine (SVM) classification and
- SVM regression.



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- loss function
- coupled with an appropriate regularization term with the help of a regularization parameter $\mu > 0$.



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The discussion below is quite brief.

Many more details can be found in specialized books or survey papers on machine learning or statistical learning such as, e.g., [EPP00, HTF09, RW06, SS02, STC04, SC08].



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The goal of the training phase of the machine learning algorithm is to determine the predictor *s* such that the empirical risk

$$R_L = \frac{1}{N} \sum_{i=1}^N L(y_i, s(\boldsymbol{x}_i))$$

is minimized.

Regularization

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Remark

The regularization term can also be interpreted as a measure of the complexity of the model (think of an eigenfunction expansion of a smooth function s with rapidly decaying eigenvalues so that high-frequency eigenfunctions contribute very little to s, i.e., s is not very complex).



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Example

The quadratic loss $L(\mathbf{y}, \mathbf{s}) = \|\mathbf{y} - \mathbf{s}\|^2$ coupled with a quadratic regularization functional leads to spline smoothing or penalized least squares.

This is also what we use for learning via RBF networks.

Regularization Theory in RKHSs

If $\mathcal{H}_{\mathcal{K}}(\Omega)$ is a RPHS with reproducing kernel \mathcal{K} we can consider $\mathbf{c}^{\mathsf{T}}\mathsf{K}\mathbf{c}$, the square of the native space norm of s, as the associated regularization term.

Theorem (Representer Theorem [KW71])

The optimal predictor s in $\mathcal{H}_{\mathcal{K}}(\Omega)$ characterized by

$$\min_{\boldsymbol{c}} \left[L(\boldsymbol{y}, \mathsf{K}\boldsymbol{c}) + \mu \boldsymbol{c}^{\mathsf{T}} \mathsf{K}\boldsymbol{c} \right],$$

can be expressed as a linear combination of kernel functions, i.e.,

$$s(\mathbf{x}) = \sum_{j=1}^{N} c_j K(\mathbf{x}, \mathbf{x}_j).$$

Here K is our usual kernel matrix and $\mathbf{y} = (y_1, \dots, y_N)^T$.

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- For general L the optimal predictor characterized in the representer theorem requires the solution of a non-trivial nonlinear optimization problem.
- If L is squared loss and we have a square system matrix K then the optimal solution is obtained by simply solving the linear system Kc = y, i.e., the empirical risk is zero and the native space norm of s is automatically minimized (see also [Fas07, Chapter 19], or the discussion below).



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Answer: Machine learning applications generally deal with data contaminated by significant errors (perhaps on the order of 10% error). This often means that the observed data look rough, as though generated by a nonsmooth function.



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Answer: Machine learning applications generally deal with data contaminated by significant errors (perhaps on the order of 10% error). This often means that the observed data look rough, as though generated by a nonsmooth function.

A common assumption is that the data were generated by a smooth function, but that the observations were corrupted by a nonsmooth error term.



This suggests we should use smooth basis functions for the approximation, but should also do some balancing to avoid fitting the nonsmooth errors.



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- Choosing μ > 0 forces c to not grow too large (and so prevents the wild oscillations which would be needed to exactly fit data from a nonsmooth function).
 - A larger μ will demand smaller c_i values and care less about fitting the data.
 - A smaller μ will more closely fit the observed data at the cost of an approximation which is more susceptible to errors in the observations, i.e., overfitting.



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Common choices of basis functions include the Gaussian, trigonometric functions and sigmoids, which have an "on/off" behavior and are used to reflect the behavior of neurons in a human brain [HTF09, Orr96].



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We will consider only so-called single layer learning algorithms using shifts of the Gaussian kernel (often called the RBF-kernel in the learning literature).





If *M* ≤ *N* copies of *K* are used with centers at locations {*z_j*}^{*M*}_{*j*=1} then

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In the interpolation setting this is zero for $\boldsymbol{c} = K^{-1} \boldsymbol{y}$.



We commonly find *c* by solving a minimization problem with squared loss of the form

$$\begin{split} \boldsymbol{c} &= \operatorname*{argmin}_{\boldsymbol{c} \in \mathbb{R}^M} \sum_{i=1}^N \left(y_i - \sum_{j=1}^M c_j K(\boldsymbol{x}_i, \boldsymbol{z}_j) \right)^2 + \mu \sum_{j=1}^M c_j^2 \\ &= \operatorname*{argmin}_{\boldsymbol{c} \in \mathbb{R}^M} \| \boldsymbol{y} - K \boldsymbol{c} \|^2 + \mu \| \boldsymbol{c} \|^2, \end{split}$$

where we specify the kernel K and regularization parameter μ beforehand.


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Remark

Since this is a convex minimization problem, the necessary condition obtained by the standard strategy of differentiating and setting the result equal to zero is also sufficient.



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where we've used $\boldsymbol{y}^T K \boldsymbol{c} = \boldsymbol{c}^T K^T \boldsymbol{y}$ since it is just a scalar.

Thus we can solve the optimization problem by solving the linear system

$$(\mathsf{K}^{\mathsf{T}}\mathsf{K} + \mu\mathsf{I}_{\mathsf{M}})\boldsymbol{c} = \mathsf{K}^{\mathsf{T}}\boldsymbol{y},\tag{1}$$

which is guaranteed to be well-defined (i.e., the inverse exists) for any $\mu > 0$.



Picking an "optimal" μ via GCV

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Generalized cross validation (GCV) is popular in the machine learning literature [CW79, GHW79, GVM97]. For our RBF networks it can be computed as [AC10, Orr96]

$$C_{GCV} = \boldsymbol{y}^T P^2 \boldsymbol{y} \frac{N}{(\text{trace P})^2}, \qquad P = I_N - K(K^T K + \mu I_M)^{-1} K^T.$$



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Remark

C_{GCV} is not an error in the true sense, though it is related to the residual:

$$\mathbf{P} \boldsymbol{y} = \boldsymbol{y} - \mathbf{K} \boldsymbol{c}, \quad \textit{i.e.,} \quad \boldsymbol{y}^T \mathbf{P}^2 \boldsymbol{y} = \| \boldsymbol{y} - \mathbf{K} \boldsymbol{c} \|_2^2.$$

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Example: Effects of μ

We take N = 50 data sampled randomly on [-1, 1] from the function

$$f(x) = (1 - 4x + 32x^2)e^{-16x^2}$$

with added normally distributed noise with zero mean and standard deviation 0.2.

MATLAB code for this example is provided in RBFNetwork1.m.





We form an RBF network using Gaussian kernels with $\varepsilon = 8$ centered at M = 15 evenly spaced points in [-1, 1].

To smooth out the noise, we consider regularization parameters μ in $[10^{-10},10^5].$



The RMS relative error is computed at 300 evenly spaced points in [-1, 1].





- "Optimal" μ from the error plot: $\mu \approx 0.56$
- "Optimal" μ from GCV: $\mu \approx$ 0.22

The network based on a larger μ oscillates less than the one produced with smaller $\mu.$

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- There is a complicated relationship between μ and ε. For some kernels μ may improve the accuracy through smoothing out noise and for others it may improve the quality by reducing ill-conditioning.
- Although μ affects the quality of the prediction s, it does not change the native space H_K (it pushes the interpolant toward a set of functions which try to minimize ||c||₂ rather than ||s||_{H_K}). In contrast, changing ε changes K and, therefore, H_K.

Example: Effects of stable basis

We use the same data as before, but now "flat" Gaussian kernels with $\varepsilon = 0.01$ so that K^TK is severely ill-conditioned.

We compare three alternative approaches:

- use of standard kernel basis { $K(\cdot, z_i)$ }, z_i evenly spaced in [-1, 1],
- use of stable basis $\{\psi_j(\cdot)\}, z_j$ evenly spaced in [-1, 1],
- use of eigenfunctions $\{\varphi_n(\cdot)\}$.

MATLAB code for this example is provided in RBFNetwork2.m.





Error compared across various regularization values μ .



MATH 590



Predictions for "optimal" μ values.



MATH 590

- The stable basis outperforms the standard basis for most μ values.
- Eigenfunctions and the stable basis essentially overlap.
- For large μ values, the prediction is dominated by the regularization component (cf. the overlap near $\mu \approx 10^4$).
- The similarity between the stable basis and eigenfunctions is due to the fact that the correction term in the HS-SVD decreases in magnitude as ε decreases.
- Even a tiny value of μ has a remarkable stabilization effect in this example.



Example: Combined effects of μ and ε

- Again, same data as before, but now we look at the effect of allowing both ε and μ to change.
- Only the eigenfunction basis is considered.
- The errors for $\mu \in [10^{-10}, 10^5]$ and $\varepsilon \in [10^{-2}, 10^1]$, along with GCV plots are shown below.
- MATLAB code for this example is provided in RBFNetwork3.m.





An optimal error exists at $\mu \approx 0.011$ and $\varepsilon \approx 2.8$, but is not predicted by GCV ("optimal" $\mu \approx 0.0017$, $\varepsilon \approx 3.7$).



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We now discuss the two main applications of support vector machines (SVMs) in the context of supervised machine learning:

- classification and
- regression.

Both of these applications can be formulated within the regularization framework outlined at the beginning of this chapter.



Standard (binary) classification

Given: a set of training data $\{(\mathbf{x}_i, \mathbf{y}_i) : i = 1, ..., N\}$ with

- measurements $\boldsymbol{x}_i \in \mathbb{R}^d$ and
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- measurements $\boldsymbol{x}_i \in \mathbb{R}^d$ and
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Find: a predictor *s* that will allow us to assign an appropriate label, either -1 or +1, to a future measurement **x**.



Example

A predictor might be

 $s(\mathbf{x}) = \operatorname{sign}(h(\mathbf{x})),$

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A typical loss function is given by the hinge loss (or soft margin loss)

$$L(y, h(\boldsymbol{x})) = \max(1 - yh(\boldsymbol{x}), 0)$$

since

$$L(y, h(\boldsymbol{x})) = 0 \quad \Longleftrightarrow \quad yh(\boldsymbol{x}) \ge 1,$$

i.e., *y* and $h(\mathbf{x})$ have the same sign and $|h(\mathbf{x})| \ge 1$ so that we have enough confidence in our prediction (see, e.g., [SS02, Chapter 3]).

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i.e., *y* and h(x) have the same sign and $|h(x)| \ge 1$ so that we have enough confidence in our prediction (see, e.g., [SS02, Chapter 3]).

An appropriate regularization term will be given by some norm of *h* (see below for more details).

Regression

We estimate continuous numeric values as discussed in the previous section.

As for RBF networks, we can use the squared loss

$$L(y, s(\mathbf{x})) = (y - s(\mathbf{x}))^2.$$



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Remark

According to the ϵ -insensitive loss function, deviations of the predicted value $s(\mathbf{x})$ from the correct value y are only penalized if they exceed ϵ , and therefore it will be possible to obtain sparse representations using only a subset of the data referred to as support vectors (more details below).

Linear Classification

The simplest predictor is given by

 $s(\mathbf{x}) = \operatorname{sign}(h(\mathbf{x}))$,

where h denotes a hyperplane — directly in input space — of the form

$$h(\boldsymbol{x}) = \boldsymbol{x}^T \boldsymbol{w} + b = 0, \qquad \boldsymbol{x} \in \mathbb{R}^d,$$

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that separates the measurements with label -1 from those with a +1.

The weights \boldsymbol{w} (which serve as the unit normal vector to the hyperplane) and the bias *b* can be determined by maximizing the margin or gap to both sides of this hyperplane (see, e.g., [HTF09, Chapter 12]).



Unconstrained minimization

Since the size of the margin is $\frac{1}{\|\boldsymbol{w}\|}$, and we want to maximize this margin, a natural regularization functional is:

minimize $\|\boldsymbol{w}\|$ (norm of the coefficients of *h*).



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Using the hinge loss function and $h(\mathbf{x}) = \mathbf{x}^T \mathbf{w} + b$, we get the unconstrained minimization problem

$$\min_{\boldsymbol{w},b}\left[\frac{1}{N}\sum_{i=1}^{N}\max\left(1-y_{i}h(\boldsymbol{x}_{i}),0\right)+\mu\frac{1}{2}\boldsymbol{w}^{T}\boldsymbol{w}\right],$$

where μ is an appropriately chosen regularization parameter.



Constrained optimization

The following constrained optimization with slack variables ξ_i is more common since it also allows us to deal with the case where the given measurements are not perfectly separable by *h*:

$$\min_{\boldsymbol{w}, b, \boldsymbol{\xi}} \left[\frac{1}{2} \boldsymbol{w}^T \boldsymbol{w} + C \sum_{i=1}^{N} \xi_i \right]$$

subject to $y_i h(\boldsymbol{x}_i) \ge 1 - \xi_i, \quad i = 1, \dots, N,$
 $\xi_i \ge 0,$

where the regularization parameter μ is transformed into $C = \frac{1}{N\mu}$.


Constrained optimization

The following constrained optimization with slack variables ξ_i is more common since it also allows us to deal with the case where the given measurements are not perfectly separable by *h*:

$$\min_{\boldsymbol{w}, b, \boldsymbol{\xi}} \left[\frac{1}{2} \boldsymbol{w}^T \boldsymbol{w} + C \sum_{i=1}^{N} \xi_i \right]$$

subject to $y_i h(\boldsymbol{x}_i) \ge 1 - \xi_i, \quad i = 1, \dots, N,$
 $\xi_i \ge 0,$

where the regularization parameter μ is transformed into $C = \frac{1}{N\mu}$.

Remark

This formulation is known in the SVM literature as the primal problem (and — ironically — as the dual problem in the optimization literature).

SVM dual problem

The SVM dual problem can be derived via Lagrange multipliers α_i (see, e.g., [HTF09, Chapter 12]) and is of the form

$$\max_{\alpha} \left(\sum_{i=1}^{N} \alpha_i - \frac{1}{2} \sum_{i=1}^{N} \sum_{j=1}^{N} \alpha_i \alpha_j \mathbf{y}_j \mathbf{y}_j \mathbf{x}_i^T \mathbf{x}_j \right)$$

subject to $\sum_{i=1}^{N} \alpha_i \mathbf{y}_i = \mathbf{0},$
 $\mathbf{0} \le \alpha_i \le C,$

where *C* is known as a box constraint and $\boldsymbol{w} = \sum_{i=1}^{N} \alpha_i y_i \boldsymbol{x}_i$ (which follows from setting the \boldsymbol{w} -gradient of the primal Lagrange multiplier functional equal to zero).

The bias *b* is given by $b = y_i - \mathbf{x}_i^T \mathbf{w}$ for any *i* such that the optimal $\alpha_i > 0$.



- For stability purposes we compute the bias by considering all qualifying indices and find b using the mean.
- The box constraint C is a free parameter which needs to be either set by the user or determined by an additional parameter optimization methods such as cross validation.



Kernel classification

Feature maps (see Chapter 2) allow us to view kernel values $K(\mathbf{x}, \mathbf{z})$ as the dot product of the transformed data in feature space, i.e., given \mathbf{x} and \mathbf{z} in input space and a feature map Φ we have

$$K(\boldsymbol{x}, \boldsymbol{z}) = \Phi(\boldsymbol{x})^T \Phi(\boldsymbol{z}).$$

Since the objective function the SVM dual problem is expressed in terms of dot products in input space we can now use the concept of feature maps and related kernels to talk about separating hyperplanes in feature space.



• The feature space is potentially infinite-dimensional (as, e.g., in the case of the Gaussian kernel) and therefore offers much more flexibility for separating the data than the finite-dimensional input space.



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- Cover's theorem [Cov65] provides a theoretical foundation for this. It ensures that data which can not be separated by a hyperplane in input space most likely will be linearly separable after being transformed to feature space by a suitable feature map.



- The feature space is potentially infinite-dimensional (as, e.g., in the case of the Gaussian kernel) and therefore offers much more flexibility for separating the data than the finite-dimensional input space.
- Cover's theorem [Cov65] provides a theoretical foundation for this. It ensures that data which can not be separated by a hyperplane in input space most likely will be linearly separable after being transformed to feature space by a suitable feature map.
- Thus, support vector machines and kernel machines, in particular — are a good tool to use in order to tackle difficult data classification problems.



Algorithms for kernel classification are essentially the same as before; simply replace the measurements x_i in input space by their transformation $\Phi(x_i)$ into feature space.



Algorithms for kernel classification are essentially the same as before; simply replace the measurements \boldsymbol{x}_i in input space by their transformation $\Phi(\boldsymbol{x}_i)$ into feature space. The separating hyperplane now is

$$h(\mathbf{x}) = \Phi(\mathbf{x})^T \mathbf{w} + b = 0, \qquad \mathbf{x} \in \mathbb{R}^d,$$

and the SVM dual problem using the transformed input data is given by

$$\max_{\alpha} \left(\sum_{i=1}^{N} \alpha_i - \frac{1}{2} \sum_{i=1}^{N} \sum_{j=1}^{N} \alpha_i \alpha_j y_i y_j \Phi(\mathbf{x}_i)^T \Phi(\mathbf{x}_j) \right)$$

subject to $\sum_{i=1}^{N} \alpha_i y_i = 0$,
 $0 \le \alpha_i \le C$.

Since we have the kernel decomposition $K(\mathbf{x}, \mathbf{z}) = \Phi(\mathbf{x})^T \Phi(\mathbf{z})$ we don't have to compute (possibly infinite) dot products in feature space, but instead just fill the kernel matrix and solve

$$\max_{\alpha} \left(\sum_{i=1}^{N} \alpha_i - \frac{1}{2} \sum_{i=1}^{N} \sum_{j=1}^{N} \alpha_i \alpha_j y_i y_j \mathcal{K}(\mathbf{x}_i, \mathbf{x}_j) \right)$$
(2)
subject to
$$\sum_{i=1}^{N} \alpha_i y_i = \mathbf{0},$$
$$\mathbf{0} \le \alpha_i \le C,$$

where, as before, *C* is the box constraint (which can be viewed as a tuning parameter) and $\boldsymbol{w} = \sum_{i=1}^{N} \alpha_i y_i \Phi(\boldsymbol{x}_i)$.



$$s(\mathbf{x}) = \operatorname{sign}(h(\mathbf{x}))$$



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= $\operatorname{sign}\left(\Phi(\mathbf{x})^{T}\mathbf{w} + b\right)$



$$\begin{aligned} s(\boldsymbol{x}) &= \operatorname{sign} \left(h(\boldsymbol{x}) \right) \\ &= \operatorname{sign} \left(\Phi(\boldsymbol{x})^T \boldsymbol{w} + b \right) \\ &= \operatorname{sign} \left(\Phi(\boldsymbol{x})^T \sum_{j=1}^N \alpha_j y_j \Phi(\boldsymbol{x}_j) + b \right) \end{aligned}$$



S(

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where *b* is obtained as before, i.e., $b = y_i - \sum_{j=1}^{N} \alpha_j y_j K(\mathbf{x}_i, \mathbf{x}_j)$ with *i* denoting the index of an α_i which is strictly between 0 and *C*. For stability purposes we can again average over all such candidates.



What does the separating hyperplane in this case look like?



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• The hyperplane will be linear only in feature space (which we usually have no concrete knowledge of). In the input space the data will be separated by a nonlinear manifold.



What does the separating hyperplane in this case look like?

- The hyperplane will be linear only in feature space (which we usually have no concrete knowledge of). In the input space the data will be separated by a nonlinear manifold.
- The representation of this manifold is sparse in the sense that not all basis functions are needed to specify it.
 - Only those centers *x_j* whose corresponding α_j are nonzero define meaningful basis functions.
 - These special centers are referred to as support vectors.



• Since the decision boundary can be expressed in terms of a limited number of support vectors, i.e., it has a sparse representation, learning is possible in very high-dimensional input spaces [SC08].

SVMs are

- robust against several types of model violations and outliers,
- computationally efficient, e.g., by using sequential minimal optimization (SMO) [Pla99] to perform the quadratic optimization task required for classification as well as regression.
- Another way to make SVMs perform more efficiently is to consider a low-rank representation for the kernel [FS02]. Below we test our own version based on eigenfunctions.



For positive definite kernels one can formulate the separating hyperplane without the bias term b. In that case the equality constraint ∑^N_{i=1} α_iy_i = 0 (which may be somewhat of a nuisance during the optimization process) can be omitted [PMR⁺01].



- For positive definite kernels one can formulate the separating hyperplane without the bias term b. In that case the equality constraint ∑^N_{i=1} α_iy_i = 0 (which may be somewhat of a nuisance during the optimization process) can be omitted [PMR⁺01].
- The primal and dual formulations each have their advantages.
 - The primal formulation (in input space) is good for large amounts of rather low-dimensional data.
 - The dual formulation (with kernels in feature space) is good for high-dimensional data (since only the number of support vectors matter).



Kernel classification

Remark

- In our numerical experiments we use Gaussian kernels.
- Linear SVM uses the dot product kernel $K(\mathbf{x}, \mathbf{z}) = \mathbf{x}^T \mathbf{z}$.
- Other popular kernels are
 - polynomial kernels of degree β in the form $K(\mathbf{x}, \mathbf{z}) = (1 + \mathbf{x}^T \mathbf{z})^{\beta}$,
 - the sigmoid kernel (or multilayer perceptron) $K(\mathbf{x}, \mathbf{z}) = \tanh(1 + \varepsilon \mathbf{x}^T \mathbf{z}).$

 Kernels may be defined via the feature map (instead of in closed form), and this feature map can be picked depending on the specific application (e.g., as a string kernel for text mining).



Outline



- 2 Radial Basis Function Networks
- 3 Classification with Support Vector Machines Theory

Olassification with Support Vector Machines — Practice

5 Support Vector Regression



Example: Simple kernel classification

This example is from [HTF09, Section 2.3] (see also help for SVMs in MATLAB's Statistics Toolbox).

We attempt to learn/classify data coming from two different populations:

- population 1, normally distributed with center at (1,0) (filled red circle) and identity covariance
- population 2, normally distributed with center at (0, 1) (filled green square) and identity covariance

Use Gaussian kernels with varying shape parameter ε and box constraint *C*.

MATLAB code for this example is provided in SVM1.m which uses SVM_Setup.m and gqr_fitsvm.m.



100 training $(\times, +)$ and 10 test points (\bigcirc, \Box)





Separating hyperplane, C = 1, $\varepsilon = 1$





Separating hyperplane, C = 1, $\varepsilon = 5$





Separating hyperplane, C = 1, $\varepsilon = 0.2$





Separating hyperplane, C = 10000, $\varepsilon = 5$





Separating hyperplane, C = 0.01, $\varepsilon = 5$





- For a fixed C, larger ε produces a more localized/detailed separator
- For a fixed ε, larger C produces a more localized/detailed separator



Effects of margin and number of support vectors

We compare the number of missed classifications (out of 20 total tests) with the margin $1/||\mathbf{w}||$ and the required number of support vectors.

We look at three experiments:

- fix C = 10000 and $\varepsilon = 0.01$,
- fix C = .6 and vary ε ,
- fix $\varepsilon = 1$ and vary *C*.

MATLAB code for this example is provided in SVM2.m which uses SVM_Setup.m and gqr_fitsvm.m.





Support vectors marked with \circ , misclassifications with \bigcirc , \Box .



MATH 590



C = 0.6, variable ε



- The margin does not appear to be useful in determining optimal ε value (it grows unboundedly as ε → 0).
- Minimizing the number of support vectors seems to suggest an optimal region for ε and it helps with computational efficiency.



 $\varepsilon = 1$, variable *C*



- All C values produce decent results.
- Large *C* values require fewer support vectors for evaluation (but more time in the optimization solution because a larger search space is used).



Use of 10-fold CV to estimate C and ε



MATLAB code for this example is provided in SVM3.m which uses SVM_Setup.m, gqr_svmcv.m and gqr_fitsvm.m.



The example just discussed uses a linearly separable pattern since the population centers (0, 1) and (1, 0) are linearly separable.

Because the $\varepsilon \to 0$ limit of Gaussians is a polynomial, it is reasonable to conclude that, with infinitely much data drawn from those populations, the optimal SVM would have $\varepsilon \to 0$ to produce a line.


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Because the $\varepsilon \to 0$ limit of Gaussians is a polynomial, it is reasonable to conclude that, with infinitely much data drawn from those populations, the optimal SVM would have $\varepsilon \to 0$ to produce a line.

We therefore next consider a different pattern which is not linearly separable.



Pattern that is not linearly separable

We want to classify data as coming from one of two populations:

- population 1 (denoted by \bigcirc and \times) with centers at $\{(0,0), (1,1), (2,0)\}$ (filled \bigcirc), and
- population 2 (denoted by □ and +), with centers at {(0,1), (1,0), (2,1)} (filled □).

Test points (large \times , +) and training points (small \times , +) are shown in the figure.

MATLAB code for this example is provided in SVM4.m which uses SVM_Setup.m, gqr_svmcv.m and gqr_fitsvm.m.





• Note that small ε can no longer produce the optimal CV residual.

 Smaller ε causes an increase in the CV residual, likely because the tendency towards polynomial behavior as ε → 0 is not desirable when learning this pattern.



SVMs are generally more popular than RBF networks.

• On the one hand, SVMs may require many fewer kernel centers for evaluation, i.e., they have a spare representation (only the nonzero coefficients must be included).



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We now

- look at that cost as a function of ε and C, and
- present a strategy for exploiting the low rank eigenfunction representation for small ε to decrease cost.



- We use the quadprog solver with the algorithm interior-point-convex from MATLAB's Optimization Toolbox with initial guess C/2 times a vector of ones.
- As always for iterative solvers, a good initial guess helps speed up convergence.



Dependence of training cost of ε and *C*

We consider the linearly separable example from above, but now use 400 training points.



MATLAB code for this example is provided in SVM4.m which uses gqr_fitsvm.m.



The solution time for the quadratic program clearly depends on ε and C.

- Very large ε and very small C seem to be solved quickly:
 - large ε because the alertkernel is very localized,
 - and small C because the solution domain is very small and quickly searched.
- Larger values of C seem to always take longer, likely because the search space is increasing.



Low-rank approximations via kernel eigenfunctions

Earlier we used the Hilbert–Schmidt SVD to avoid ill-conditioning of the kernel matrix K.

However, this does not help here because the inverse of the kernel matrix is not needed during the quadratic program solution.



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Instead, we may use the eigenfunction expansion

 $\mathsf{K} = \Phi \Lambda \Phi^T$

to produce a low-rank approximation of K and exploit this structure to decrease the cost of the quadratic program.



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Instead, we may use the eigenfunction expansion

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to produce a low-rank approximation of K and exploit this structure to decrease the cost of the quadratic program.

Given *N* input points and a small ε , only a very low number *M* of eigenfunctions may be needed to accurately approximate K.



The quadratic program (2) can be written in matrix form as

$$\begin{split} \min_{\boldsymbol{\alpha}} & \frac{1}{2} \boldsymbol{\alpha}^T \mathsf{D}_{\boldsymbol{y}} \mathsf{K} \mathsf{D}_{\boldsymbol{y}} \boldsymbol{\alpha} - \boldsymbol{e}^T \boldsymbol{\alpha} \\ \text{subject to } & \boldsymbol{y}^T \boldsymbol{\alpha} = \mathbf{0}, \\ & \boldsymbol{\alpha} \in [\mathbf{0}, \boldsymbol{C}]^N, \end{split}$$
(3)

where D_y is a diagonal matrix with y on the diagonal, and e is a vector of all ones.



The quadratic program (2) can be written in matrix form as

$$\begin{split} \min_{\alpha} \frac{1}{2} \alpha^{T} \mathsf{D}_{\boldsymbol{y}} \mathsf{K} \mathsf{D}_{\boldsymbol{y}} \alpha - \boldsymbol{e}^{T} \alpha \qquad (3) \\ \text{subject to } \boldsymbol{y}^{T} \alpha = 0, \\ \alpha \in [0, C]^{N}, \end{split}$$

where D_y is a diagonal matrix with y on the diagonal, and e is a vector of all ones.

Using $K\approx (\Lambda^{1/2}\Phi)^T(\Lambda^{1/2}\Phi),$ we can rephrase this problem as [FS02, ZTK08]

$$\begin{split} \min_{\eta,\alpha} \frac{1}{2} \begin{pmatrix} \eta^{T} & \alpha^{T} \end{pmatrix} \begin{pmatrix} \mathsf{I}_{M} & \mathsf{0} \\ \mathsf{0} & \mathsf{0} \end{pmatrix} \begin{pmatrix} \eta \\ \alpha \end{pmatrix} - \begin{pmatrix} \mathsf{0} & \boldsymbol{e}^{T} \end{pmatrix} \begin{pmatrix} \eta \\ \alpha \end{pmatrix} \qquad (4) \\ \text{subject to } \begin{pmatrix} \mathsf{0} & \boldsymbol{y}^{T} \\ -\mathsf{I}_{M} & \Lambda^{1/2} \Phi^{T} \mathsf{D}_{\boldsymbol{y}} \end{pmatrix} \begin{pmatrix} \eta \\ \alpha \end{pmatrix} = \mathsf{0}, \\ \alpha \in [\mathsf{0}, C]^{N}, \quad \eta \in \mathbb{R}^{M}. \end{split}$$

Although this system is of size N + M (and the original system was only size N), the cost of solving this system may be much lower because of the extremely simple structure of the Hessian.

This sparsity, in comparison to H which may be fully dense, allows for cheap matrix-vector products and decompositions, both of which may all for a faster quadratic program solve.

Note that the η values are inconsequential in making predictions with the SVM.



Low-rank vs. full-rank approximation

We use the same setup as before (given in SVM5.m).

We study the cost of solving the quadratic program and training the SVM. Minimizing this cost is an important topic in machine learning (see, e.g., [YDD04, FL02, LLZ⁺11]).

Increasingly large sets of input points are considered and the cost of solving the full rank problem (3) is compared to solving the low rank problem (4).

The kernel is parameterized with $\varepsilon = .01$ and C = 1 and the eigenfunctions of the Gaussian use $\alpha = 10^6$.







Outline



- 2 Radial Basis Function Networks
- 3 Classification with Support Vector Machines Theory
- 4 Classification with Support Vector Machines Practice
- 5 Support Vector Regression



Linear support vector regression

As for classification, we again start with a linear approximation and assume that

$$s(\mathbf{x}) = \mathbf{x}^T \mathbf{w} + b.$$

If we use the ϵ -insensitive loss function

$$L(y, s(\mathbf{x})) = \max(|y - s(\mathbf{x})| - \epsilon, 0)$$

then the primal unconstrained minimization problem is given by

$$\min_{\boldsymbol{w},b} \left[\frac{1}{N} \sum_{i=1}^{N} \max\left(|\boldsymbol{y}_i - \boldsymbol{s}(\boldsymbol{x}_i)| - \epsilon, 0 \right) + \mu \frac{1}{2} \boldsymbol{w}^T \boldsymbol{w} \right],$$

where, as before, μ is an appropriately chosen regularization parameter.



Constrained minimization problem

Using slack variables as in the classification case we have the analogous constrained minimization problem

$$\min_{\boldsymbol{w}, b, \boldsymbol{\xi}, \boldsymbol{\xi}^*} \left[\frac{1}{2} \boldsymbol{w}^T \boldsymbol{w} + C \sum_{i=1}^N (\xi_i + \xi_i^*) \right]$$
subject to $\boldsymbol{s}(\boldsymbol{x}_i) - \boldsymbol{y}_i \le \epsilon + \xi_i, \quad i = 1, \dots, N,$
 $\begin{array}{l} \boldsymbol{y}_i - \boldsymbol{s}(\boldsymbol{x}_i) \le \epsilon + \xi_i^*, \quad i = 1, \dots, N, \\ \xi_i, \xi_i^* \ge 0. \end{array}$



Dual problem

In the dual formulation we need to solve the constrained quadratic programming problem

$$\min_{\boldsymbol{\alpha},\boldsymbol{\alpha}^*} \epsilon \sum_{i=1}^N (\alpha_i^* + \alpha_i) - \sum_{i=1}^N \mathbf{y}_i (\alpha_i^* - \alpha_i) + \frac{1}{2} \sum_{i=1}^N \sum_{j=1}^N (\alpha_i^* - \alpha_i) (\alpha_j^* - \alpha_j) \mathbf{x}_i^T \mathbf{x}_j$$

subject to $0 \leq \alpha_i, \alpha_i^* \leq C$,

$$\sum_{i=1}^{N} (\alpha_i^* - \alpha_i) = \mathbf{0}.$$



Dual problem

In the dual formulation we need to solve the constrained quadratic programming problem

$$\begin{split} \min_{\alpha,\alpha^*} \epsilon \sum_{i=1}^N (\alpha_i^* + \alpha_i) - \sum_{i=1}^N y_i (\alpha_i^* - \alpha_i) + \frac{1}{2} \sum_{i=1}^N \sum_{j=1}^N (\alpha_i^* - \alpha_i) (\alpha_j^* - \alpha_j) \boldsymbol{x}_i^T \boldsymbol{x}_j \\ \text{subject to } 0 \le \alpha_i, \alpha_i^* \le C, \\ \sum_{i=1}^N (\alpha_i^* - \alpha_i) = 0. \end{split}$$

Once we've found the dual variables
$$\alpha_i$$
 and α_i^* , the SVM regression function is given by

$$\boldsymbol{s}(\boldsymbol{x}) = \boldsymbol{x}^T \boldsymbol{w} + \boldsymbol{b} = \sum_{i=1}^{N} (\alpha_i^* - \alpha_i) \boldsymbol{x}^T \boldsymbol{x}_i + \boldsymbol{b},$$

i.e.,
$$\boldsymbol{w} = \sum_{i=1}^{N} (\alpha_i^* - \alpha_i) \boldsymbol{x}_i$$

i=1



• The computation of the bias term b follows from the KKT conditions (for details see [SS02]) and is similar in spirit to the classification setting, i.e.,

$$b = y_i - \boldsymbol{x}_i^T \boldsymbol{w} - \epsilon \quad \text{for } \alpha_i \in (0, C),$$

$$b = y_i - \boldsymbol{x}_i^T \boldsymbol{w} + \epsilon \quad \text{for } \alpha_i^* \in (0, C).$$

As before, any one of these will theoretically suffice, but for stability reasons it is better to compute b via an average over all candidates.

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As in the classification setting, α^{*}_i − α_i ≠ 0 only for some i, and the corresponding measurements x_i are called the support vectors.

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- As in the classification setting, α^{*}_i − α_i ≠ 0 only for some i, and the corresponding measurements x_i are called the support vectors.
- For more details see, e.g., [HTF09, Chapter 12], [SS02, Chapter 9].

Nonlinear support vector regression

As for classification, we obtain a nonlinear "kernelized" regression fit if we map the data into feature space and then use kernels.

This is straightforward and completely analogous to the classification setting.



The resulting dual problem is

$$\min_{\alpha,\alpha^*} \epsilon \sum_{i=1}^N (\alpha_i^* + \alpha_i) - \sum_{i=1}^N y_i (\alpha_i^* - \alpha_i) + \frac{1}{2} \sum_{i=1}^N \sum_{j=1}^N (\alpha_i^* - \alpha_i) (\alpha_j^* - \alpha_j) \mathcal{K}(\boldsymbol{x}_i, \boldsymbol{x}_j)$$

subject to $0 \le \alpha_i, \alpha_i^* \le C$,

$$\sum_{i=1}^{N} (\alpha_i^* - \alpha_i) = \mathbf{0},$$

so that

$$\boldsymbol{s}(\boldsymbol{x}) = \Phi(\boldsymbol{x}^T) \boldsymbol{w} + \boldsymbol{b} = \sum_{i=1}^{N} (\alpha_i^* - \alpha_i) K(\boldsymbol{x}, \boldsymbol{x}_i) + \boldsymbol{b},$$

i.e., $\boldsymbol{w} = \sum_{i=1}^{N} (\alpha_i^* - \alpha_i) \Phi(\boldsymbol{x}_i)$ and

$$b = y_i - \sum_{j=1}^{N} (\alpha_j^* - \alpha_j) K(\boldsymbol{x}_i, \boldsymbol{x}_j) - \epsilon \quad \text{for } \alpha_i \in (0, C),$$

$$b = y_i - \sum_{j=1}^{N} (\alpha_j^* - \alpha_j) K(\boldsymbol{x}_i, \boldsymbol{x}_j) + \epsilon \quad \text{for } \alpha_i^* \in (0, C).$$



MATH 590

Many more details on all aspects of machine learning can be found, e.g., in the

- books [Alp09, HTF09, RW06, SS02, STC04, SC08] or
- survey papers [EPP00, MMn06, Orr96].



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