Recap

- ▶ We are discussing the SVM method for learning classifiers.
- ▶ Objective is to learn the optimal hyperplane one that maximizes the 'margin of separation'.
- ▶ It can be formulated as a constrained optimization problem.

Linear SVM - data linearly separable

▶ The optimal hyperplane is a solution of

minimize
$$\frac{1}{2}W^TW$$
 subject to $y_i(W^TX_i+b)\geq 1, \ i=1,\ldots,n$

We solve the dual given by

$$\max_{\pmb{\mu}} \qquad q(\pmb{\mu}) = \sum_{i=1}^n \ \mu_i - \frac{1}{2} \sum_{i,j=1}^n \ \mu_i \mu_j y_i y_j X_i^T X_j$$
 subject to
$$\mu_i \geq 0, \quad i = 1, \dots, n, \quad \sum_{i=1}^n \ y_i \mu_i = 0$$

▶ Then the final solution is:

$$W^* = \sum_{i} \mu_i^* y_i X_i, \ b^* = y_j - X_j^T W^*, \ j \text{ such that } \mu_j > 0$$

Linear SVM

► The primal problem is

minimize
$$\frac{1}{2}W^TW + C\sum_{i=1}^n \xi_i$$
 subject to
$$y_i(W^TX_i + b) \geq 1 - \xi_i, \quad i = 1, \dots, n$$

$$\xi_i \geq 0, \quad i = 1, \dots, n$$

► The dual problem is:

$$\max_{\boldsymbol{\mu}} \qquad q(\boldsymbol{\mu}) = \sum_{i=1}^n \ \mu_i - \frac{1}{2} \sum_{i,j=1}^n \ \mu_i \mu_j y_i y_j X_i^T X_j$$
 subject to
$$0 \le \mu_i \le C, \quad i = 1, \dots, n, \quad \sum_i y_i \mu_i = 0$$

▶ We solve dual and the final optimal hyperplane is

$$W^* = \sum_i \mu_i^* y_i X_i,$$

$$b^* = y_i - X_i^T W^*, (j \ s.t. \ 0 < \mu_i < C).$$

The Linear SVM

▶ Given training data, $(X_i, y_i), i = 1, \dots, n$, we solve

$$\max_{\boldsymbol{\mu}} \qquad q(\boldsymbol{\mu}) = \sum_{i=1}^{n} \mu_i - \frac{1}{2} \sum_{i,j=1}^{n} \mu_i \mu_j y_i y_j X_i^T X_j$$

subject to $0 \le \mu_i \le C, \quad i = 1, \dots, n, \quad \sum y_i \mu_i = 0$

▶ The SVM is a linear classifier specified by W^*, b^* :

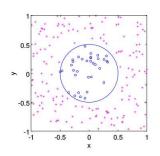
$$W^* = \sum \mu_i^* y_i X_i, \quad b^* = y_j - X_j^T W^*, \ (j \ s.t. \ 0 < \mu_j < C)$$

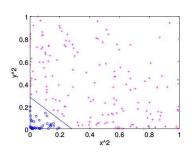
ightharpoonup Given a new pattern, X, its class is determined by sign of

$$f(X) = X^T W^* + b^* = \sum_{i} \mu_i^* y_i X_i^T X + b^*$$
$$= \sum_{i:\mu_i>0} \mu_i^* y_i X_i^T X + (y_j - \sum_{i:\mu_i>0} \mu_i^* y_i X_i^T X_j)$$

- We first formulated SVM for linearly separable case to understand the idea.
- ▶ Then, using slack variables, ξ_i , we have a linear SVM where we make no assumptions about separability.
- ▶ In the dual, the only difference is an upperbound on μ_i .
- Using this formulation, we can find 'best' hyperplane classifier.
- The next question is how can we learn non-linear classifiers?
- ightharpoonup Recall that the SVM idea is to transform X_i into some other high-dimensional space and learn a linear classifier there.

Transforming Patterns to become Linearly Separable





Non-linear classifiers

- ▶ In general, we can use a mapping, $\phi: \Re^m \to \Re^{m'}$.
- In $\Re^{m'}$, the training set is $\{(Z_i, y_i), i = 1, \dots, n\}, Z_i = \phi(X_i).$
- We can find optimal hyperplane by solving the dual (replacing $X_i^T X_j$ with $Z_i^T Z_j$).
- ▶ The dual problem now would be the following.

$$\max_{\mu} \qquad q(\mu) = \sum_{i=1}^{n} \mu_{i} - \frac{1}{2} \sum_{i,j=1}^{n} \mu_{i} \mu_{j} y_{i} y_{j} \phi(X_{i})^{T} \phi(X_{j})$$

subject to
$$0 \le \mu_i \le C, \quad i = 1, \dots, n, \quad \sum_{i=1}^n y_i \mu_i = 0$$

- ▶ The optimization is over \Re^n (with quadratic cost function & linear constraints) **irrespective of** ϕ **and** m'.
- But computationally expensive?

Kernel function

▶ Suppose we have a function, $K: \Re^m \times \Re^m \to \Re$, that satisfies

$$K(X_i, X_j) = \phi(X_i)^T \phi(X_j)$$

Called Kernel function.

- ▶ Suppose computation of $K(X_i, X_j)$ is about as expensive as that of $X_i^T X_j$.
- ▶ Replacing $\phi(X_i)^T \phi(X_j)$ by $K(X_i, X_j)$, we can solve dual without ever computing any $\phi(X_i)$. Efficient for obtaining optimal hyperplane.
- ▶ What about storing W^* ? Computing $\phi(X)^TW^*$ for new patterns?

Kernel function based classifier

- ▶ Let μ_i^* be soln of Dual. Then $W^* = \sum \mu_i^* y_i \phi(X_i)$.
- We also have

$$b^* = y_j - \phi(X_j)^T W^* = y_j - \sum_i \mu_i^* y_i \phi(X_i)^T \phi(X_j)$$

▶ Given a new pattern X we only need to compute

$$f(X) = \phi(X)^T W^* + b^*$$

$$= \sum_{i} \mu_i^* y_i \phi(X_i)^T \phi$$

$$= \sum_{i} \mu_{i}^{*} y_{i} \phi(X_{i})^{T} \phi(X) + (y_{j} - \phi(X_{j})^{T} W^{*})$$

$$= \sum_{i} \mu_{i}^{*} y_{i} \phi(X_{i})^{T} \phi(X) + \left(y_{j} - \phi(X_{j})^{T} W^{*}\right)$$

$$= \sum_{i} \mu_i^* y_i K(X_i, X) + \left(y_j - \sum_{i} \mu_i^* y_i K(X_i, X_j) \right)$$
$$= \sum_{i:\mu_i>0} \mu_i^* y_i K(X_i, X) + \left(y_j - \sum_{i:\mu_i>0} \mu_i^* y_i K(X_i, X_j) \right)$$

- ▶ This is an interesting way of learning nonlinear classifiers.
- We solve the dual whose dimension is n, number of examples. (We do not need $\phi(X_i)$ for solving dual)
- All we need to store are:
 - non-zero Lagrange multipliers: $\mu_i^* > 0$,
 - Support vectors: X_i , i s.t. $\mu_i^* > 0$.
- ▶ Then, given an X, we compute

$$f(X) = \sum_{i:\mu_i>0} \mu_i^* y_i K(X_i, X) + \left(y_j - \sum_{i:\mu_i>0} \mu_i^* y_i K(X_i, X_j) \right)$$

and classify X based on sign of f(X).

▶ Never need to enter ' $\phi(X)$ ' space!

Support Vector Machine

- ▶ Obtain μ_i^* by solving the Dual with $\phi(X_i)^T \phi(X_j)$ replaced by $K(X_i, X_j)$. (Choose a suitable Kernel function. Use 'penalty const', C as needed).
- ▶ Store non-zero μ_i^* and the corresponding support vectors.
- Classify any new pattern X by sign of

$$f(X) = \sum \mu_i^* y_i K(X_i, X) + \left(y_j - \sum_i \mu_i^* y_i K(X_i, X_j) \right)$$

- ▶ If we have a suitable Kernel function, we never need to compute $\phi(X)$.
- ▶ The range space of ϕ can even be infinite dimensional!

Example kernel function

- ▶ We start with an example kernel function in \Re^2 .
- Consider $K(X_i, X_j) = (1 + X_i^T X_j)^2$.
- ▶ Let $X_i = (x_{i1}, x_{i2})^T \in \Re^2$ and similarly for X_i .
- ► Then

$$K(X_i, X_j) = (1 + x_{i1}x_{j1} + x_{i2}x_{j2})^2$$

• We now show that there exists a mapping ϕ such that $K(X_i, X_j) = \phi(X_i)^T \phi(X_j)$.

▶ Consider $\phi: \Re^2 \to \Re^6$ given by

$$Z=\phi(X)=\begin{bmatrix} 1 & x_1^2 & x_2^2 & \sqrt{2}x_1 & \sqrt{2}x_2 & \sqrt{2}x_1x_2 \end{bmatrix}$$
 (Here, $X=(x_1 \ x_2)\in\Re^2$).

- ▶ It is easy to see that a linear discriminant function in terms of Z (i.e., in \Re^6) would be a quadratic discriminant function in terms of X (i.e., in \Re^2).
- ▶ Now we show that

$$K(X_i, X_j) = (1 + X_i^T X_j)^2 = Z_i^T Z_j = \phi(X_i)^T \phi(X_j)$$

► Recall

$$Z_i = \phi(X_i) = \begin{bmatrix} 1 & x_{i1}^2 & x_{i2}^2 & \sqrt{2}x_{i1} & \sqrt{2}x_{i2} & \sqrt{2}x_{i1}x_{i2} \end{bmatrix}$$

We have

$$Z_i^T Z_j = 1 + x_{i1}^2 x_{j1}^2 + x_{i2}^2 x_{j2}^2 + 2x_{i1} x_{j1} + 2x_{i2} x_{j2} + 2x_{i1} x_{i2} x_{j1} x_{j2}$$

$$= (1 + x_{i1} x_{j1} + x_{i2} x_{j2})^2$$

$$= (1 + X_i^T X_j)^2 = K(X_i, X_j)$$

- ▶ Easy to see it works for $X \in \Re^n$ in general.
- ► Thus $K(X_i, X_j) = (1 + X_i^T X_j)^2$ results in a quadratic discriminant function or a quadratic classifier.

- From this example, it is also easy to see that for a given Kernel function, the mapping ϕ (or the dimension of its range space) is not unique.
- ▶ Consider the same Kernel fn $K(X_i, X_j) = (1 + X_i^T X_j)^2$.
- Consider the mapping $\phi: \Re^2 \to \Re^7$ given by

$$Z = \phi(X) = \begin{bmatrix} 1 & \sqrt{2}x_1 & \sqrt{2}x_2 & x_1^2 & x_2^2 & x_1x_2 & x_1x_2 \end{bmatrix}$$

It is easy to see that this mapping also works.

- ▶ We saw that the Kernel $K(X, X') = (1 + X^T X')^2$ results in a quadratic discriminant function (in the original feature space)
- ▶ This is because the effective ϕ function is such that each x_ix_j term is a component of $\phi(X)$.
- ▶ Thus, if $X \in \Re^m$, then any reasonable ϕ function corresponding to this kernel would have range space with dimension $O(m^2)$.
- ▶ Hence, $\phi(X_i)^T \phi(X_i)$ would need $O(m^2)$ multiplications.
- ▶ If we are using a linear SVM, we only need $X_i^T X_j$ which needs m multiplications.
- When we use the Kernel for the quadratic case, we need only m+1 multiplications.

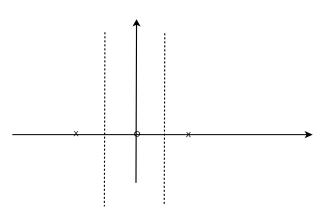
Example

- We will first consider a very simple example problem in \Re^2 to get a feel for the method of obtaining SVM.
- Suppose we have 3 examples:

$$X_1 = (-1,0), \quad X_2 = (1,0), \quad X_3 = (0,0)$$

with $y_1 = y_2 = +1$ and $y_3 = -1$.

Example



- ▶ As is easy to see, a linear classifier is not sufficient here.
- Suppose we use the Kernel function: $K(X, X') = (1 + X^T X')^2$.

Recall, the examples are

$$X_1 = (-1,0), \quad X_2 = (1,0), \quad X_3 = (0,0)$$

▶ The objective function involves $K(X_i, X_j)$. These are given in a matrix below.

$$\left[(1 + X_i^T X_j)^2 \right] = \begin{bmatrix} 4 & 0 & 1 \\ 0 & 4 & 1 \\ 1 & 1 & 1 \end{bmatrix}$$

▶ The objective function to be maximized is

$$q(\boldsymbol{\mu}) = \sum_{i=1}^{3} \mu_i - \frac{1}{2} \sum_{i,j=1}^{3} \mu_i \mu_j y_i y_j K(X_i, X_j)$$
$$= \sum_{i=1}^{3} \mu_i - \frac{1}{2} (4\mu_1^2 + 4\mu_2^2 + \mu_3^2 - 2\mu_1 \mu_3 - 2\mu_2 \mu_3)$$

▶ The constraints are

$$\mu_1 + \mu_2 - \mu_3 = 0;$$
 and $-\mu_i \le 0, i = 1, 2, 3.$

▶ The lagrangian for this problem is

$$L(\boldsymbol{\mu}, \lambda, \boldsymbol{\alpha}) = q(\boldsymbol{\mu}) + \lambda(\mu_1 + \mu_2 - \mu_3) - \sum_{i=1}^{3} \alpha_i \mu_i$$

Here, λ is the Lagrange multiplier for the equality constraint and α_i are the langrange multipliers for the inequality constraints.

- ▶ Using Kuhn-Tucker conditions, we have $\frac{\partial L}{\partial \mu_i} = 0$ and $\mu_1 + \mu_2 \mu_3 = 0$.
- ▶ This gives us four equations; we have 7 unknowns (Three μ_i , three α_i and λ).
- ▶ By complementary slackness, we have $\alpha_i \mu_i = 0$. Essentially, we need to guess which $\mu_i > 0$.

- ▶ In this simple problem we know all $\mu_i > 0$.
- \triangleright This is because all X_i would be support vectors.
- ▶ Hence we take all $\alpha_i = 0$.
- ▶ We have now four unknowns: $\mu_1, \mu_2, \mu_3, \lambda$.
- ▶ Using $\frac{\partial L}{\partial \mu_i} = 0$, i = 1, 2, 3 and feasibility, we can solve for μ_i .

Recall

$$L(\boldsymbol{\mu}, \lambda, \boldsymbol{\alpha}) = q(\boldsymbol{\mu}) + \lambda(\mu_1 + \mu_2 - \mu_3) - \sum_{i=1}^{3} \alpha_i \mu_i$$

$$q(\boldsymbol{\mu}) = \sum_{i=1}^{3} \mu_i - \frac{1}{2} (4\mu_1^2 + 4\mu_2^2 + \mu_3^2 - 2\mu_1\mu_3 - 2\mu_2\mu_3)$$

 $1 - \mu_3 + \mu_1 + \mu_2 - \lambda = 0$

 $1 - 4\mu_1 + \mu_3 + \lambda = 0$ $1 - 4\mu_2 + \mu_3 + \lambda = 0$

Now $\frac{\partial L}{\partial u} = 0$, i = 1, 2, 3 and feasibility give

$$\mu_1+\mu_2-\mu_3 \ = \ 0$$

 These give us $\lambda=1$ and $\mu_3=2\mu_1=2\mu_2.$

- ▶ Thus we get $\mu_1 = \mu_2 = 1$ and $\mu_3 = 2$.
- ► This completely determines the SVM

- ▶ If we used the penalty constant with $C \ge 2$ we get the same solution.
 - (If C < 2, we can not get this solution).
- ► The calssification of any X by this SVM is by the sign of f(X):

$$f(X) = \sum_{i} \mu_{i} y_{i} K(X_{i}, X) + b^{*}$$
$$= K(X_{1}, X) + K(X_{2}, X) - 2K(X_{3}, X) + b^{*}$$

▶ Let us first calculate b*.

Recall the formula

$$b^* = y_j - \sum_i \mu_i y_i K(X_i, X_j), \quad j \ s.t. \ 0 < \mu_j$$

Recall

$$[K(X_i, X_j)] = [(1 + X_i^T X_j)^2] = \begin{bmatrix} 4 & 0 & 1 \\ 0 & 4 & 1 \\ 1 & 1 & 1 \end{bmatrix}$$

- With j = 1 we get $b^* = 1 (4 + 0 2) = -1$.
- ▶ With j = 3 we get $b^* = -1 (1 + 1 2) = -1$.
- ▶ If we solved our optimization problem correctly, we should get same b*!

- ▶ We have $X_1 = (-1,0), X_2 = (1,0), X_3 = (0,0)$ and $K(X,X') = (1+X^TX')^2$.
- Hence, taking $X = (x_1, x_2)^T$, we have

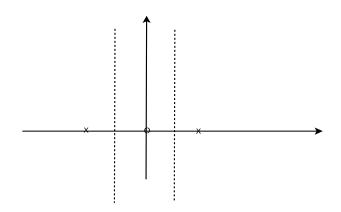
$$f(X) = K(X_1, X) + K(X_2, X) - 2K(X_3, X) + b^*$$

= $(1 - x_1)^2 + (1 + x_1)^2 - 2(1) - 1$
= $2x_1^2 - 1$

lacksquare Hence this SVM will assign class +1 to $X=(x_1,x_2)^T$ if

$$2x_1^2 \ge 1$$
 or $|x_1| \ge \frac{1}{\sqrt{2}}$

Recall Example Patterns



▶ Hence this SVM will assign class +1 to $X = (x_1, x_2)^T$ if

$$2x_1^2 \ge 1$$
 or $|x_1| \ge \frac{1}{\sqrt{2}}$

- Why not $|x_1| \ge (1/2)$?
- We are maximizing margin of the hyperplane in x^2 '-space.
- ► The final SVM is intuitively very reasonable and we solve essentially the same problem whether we are seeking a linear classifier or a nonlinear classifier.

Kernel functions

- ▶ How do we obtain Kernel functions in general?
- What kind of symmetric functions capture the inner product in an appropriate space?
- We look at two important charectarizations for Kernel functions.

Mercer Kernels

Mercer Theorem:

Given a symmetric function, $K: \Re^m \times \Re^m \to \Re$, there exists an inner product space $\mathcal H$ and a mapping $\phi: \Re^m \to \mathcal H$ so that

$$K(X_1, X_2) = \phi(X_1)^T \phi(X_2)$$

if for all square-integrable functions g,

$$\int K(X_1, X_2)g(X_1)g(X_2)dX_1 dX_2 \ge 0.$$

Positive definite kernels

- ▶ Let \bar{K} be a $n \times n$ matrix with $\bar{K}_{i,j} = K(X_i, X_j)$.
- ▶ A **positive definite kernel** is the function K such that \bar{K} is positive semi-definite for all n and all data sets $\{X_1, \ldots, X_n\}$.
- ▶ That is, given any n, and any feature vectors, X_1, \dots, X_n , we have, for all scalars c_1, \dots, c_n ,

$$\sum_{i,j=1}^{n} c_i c_j K(X_i, X_j) \ge 0$$

▶ If input space is compact, both these notions are same.

- ▶ Now we use Mercer's theorem to show that the function we gave earlier would be a Kernel function.
- Consider the function

$$K(U, V) = (U^T V)^p = \left(\sum_{i=1}^m u_i v_i\right)^p$$

where p > 0 is an integer and $U = [u_1 \cdots u_m]^T$ and $V = [v_1 \cdots v_m]^T$ are in \Re^m .

▶ We want to show that this satisfies the Mercer theorem.

• By expanding the $(U^TV)^p$ we get an expression

$$\left(\sum_{i=1}^{m} u_i v_i\right)^p = \sum_{r_1, \dots, r_m} \frac{p!}{r_1! r_2! \cdots r_m!} \prod_{i=1}^{m} (u_i v_i)^{r_i}$$

where the summation is over all non-negative integers, $r1, \cdots, rm$ such that

$$r1 + r2 + \cdots + rm = p$$

We need to show

$$\int_{\Re^m} \int_{\Re^m} \left(\sum_{i=1}^m u_i v_i \right)^p g(U) g(V) dU dV > 0.$$

- ▶ This becomes a sum of integrals by expanding $(\sum u_i v_i)^p$.
- A typical term here is

$$\frac{p!}{r_1! \, r_2! \cdots r_m!} \int \int (u_1 v_1)^{r_1} (u_2 v_2)^{r_2} \cdots (u_m v_m)^{r_m} g(U) g(V) dU dV
= \frac{p!}{r_1! \, r_2! \cdots r_m!} \int (u_1)^{r_1} (u_2)^{r_2} \cdots (u_m)^{r_m} g(U) dU
\int (v_1)^{r_1} (v_2)^{r_2} \cdots (v_m)^{r_m} g(V) dV
= \frac{p!}{r_1! \, r_2! \cdots r_m!} \left(\int u_1^{r_1} u_2^{r_2} \cdots u_m^{r_m} g(U) dU \right)^2 \ge 0$$

Now consider the function

$$K(U, V) = \sum_{j=0}^{p} a_j (U^T V)^j, \ a_j \ge 0$$

We can show this also satisfies Mercer theorem

$$\int \sum_{j=0}^{p} a_{j} (U^{T}V)^{j} g(U) g(V) dU dV
= \sum_{j=0}^{p} a_{j} \int (U^{T}V)^{j} g(U) g(V) dU dV
\ge 0$$

Hence functions of the form

$$K(X_1, X_2) = \sum_{j=0}^{p} a_j (X_1^T X_2)^j, \ a_j \ge 0$$

are kernels (satisfying Mercer's theorem).

▶ A special case is

$$K(X_1, X_2) = (1 + X_1^T X_2)^p$$

which is an example we considered earlier.

This is called a polynomial kernel.

We showed that the function

$$K(X_1, X_2) = (1 + X_1^T X_2)^p$$

satisfies the Mercer theorem.

- Hence it is a (mercer) kernel.
- It is easy to show that it satisfies the definition of a positive definite kernel
- Hence it is also a positive definite kernel

Now consider the functions of the type

$$K(U, V) = \sum_{j=0}^{\infty} a_j (U^T V)^j, \ a_j \ge 0$$

- Our proof only involved interchanging integration and summation.
- For finite sum it is always possible.
- ► For infinite sum, a sufficient condition is that the above sum is uniformly convergent
- ▶ Then the above would also satisfy Mercer's theorem.

Consider the function

$$K(X_1, X_2) = e^{-\frac{(X_1 - X_2)^T (X_1 - X_2)}{2\sigma^2}}$$

▶ We can show it satisfies the theorem by noting

$$e^{-(X_1-X_2)^T(X_1-X_2)} = e^{-X_1^TX_1} e^{-X_2^TX_2} e^{2X_1^TX_2},$$

and

$$e^{2X_1^T X_2} = \sum_{p=0}^{\infty} \frac{(2X_1^T X_2)^p}{p!}$$

Some Popular Kernel functions

Polynomial kernel:

$$K_p(X_1, X_2) = (1 + X_1^T X_2)^p$$

Gaussian kernel

$$K_G(X_1, X_2) = e^{-\frac{||X_1 - X_2||^2}{\sigma^2}}$$

Generalization abilities

- SVM idea: learn linear classifier in a transformed space.
- ▶ We said that naively transforming patterns into a high dimensional space does not work.
- There were two issues in that 'curse of dimensionality'
- Computational complexity 'elegantly' solved by the kernel trick.
- But, does SVM generalize well?
- ▶ We are finding a hyperplane in a very high dimensional space. Do we need very large number of examples?

- ▶ In practice, SVMs perform well.
- They can learn classifiers that do well on test data without needing (correspondingly) large number of examples.
- ► The reason, essentially, is that we learn a hyperplane with large margin.
- ▶ There are different ways to analyze this.
- We would just state one theoretical result.

Some theoretical results

- Let P_{err}^n be the error rate on a test set for an SVM trained with n random examples.
- ► Then we can show (for SVM with no slack variables)

$$EP_{\mathsf{err}}^n \le \min\left(\frac{s}{n}, \frac{[R^2||W||^2]}{n}, \frac{m}{n}\right)$$

where s is number of support vectors, R is the radius of smallest sphere enclosing all examples, $||W||^{-2}$ is the margin of the maximum margin hyperplane (in the feature space of dimension m).

We have

$$EP_{\mathsf{err}}^n \le \min\left(\frac{s}{n}, \frac{[R^2||W||^2]}{n}, \frac{m}{n}\right)$$

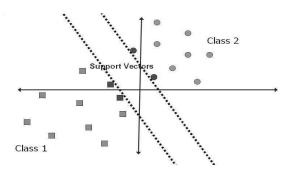
- Optimal hyperplane may generalize well because
 - ▶ 'data compression' is large
 - Margin is large
 - dimension of feature space is small

- ► The fraction of support vectors is an upperbound on expected generalization error.
- ▶ This is useful as a confidence measure on the learnt SVM.
- ► This essentially comes from the concept of 'stability' of a learning algorithm.

- Let $\{X_1, \dots, X_n\}$ be a data set from which we learnt a classifier.
- ▶ let Z be another random example. If we gave $\{Z, X_2, \cdots, X_n\}$ as the training set we expect to learn a very similar classifier.
- ▶ Similarly we can make n data sets by replacing each X_i in turn with Z.
- ▶ The learning algorithm is stable if the errors on the set $\{Z, X_1, \dots, X_n\}$ of all these classifiers are close.
- One can put a bound on generalization error based on the difference of errors of these classifiers.

Optimal hyperplane and the support Vectors

The optimal hyperplane does not change if we remove any 'non-support-vector' from the training data.

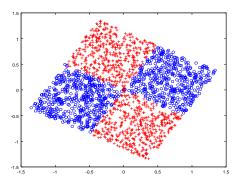


How good is SVM method?

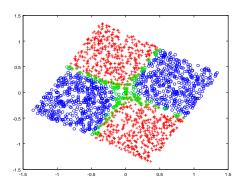
- A very competitive method for tackling many PR problems.
- ► Learning a nonlinear classifier only involves choosing a Kernel function.
- ▶ User needs to make choice of parameters kernel function and C. Also parametrs of the Opt technique. Bad choices result in 'overfitting'.
- Support vectors are an imporatnt extra information we get.

An example

▶ Consider the training set for a 2-class problem in \Re^2 as below.

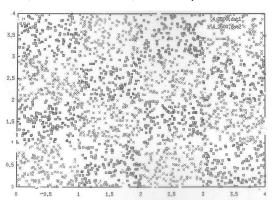


► The support vectors are shown below. (Gaussian Kernel)

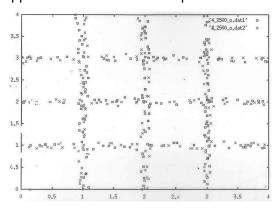


Another example

▶ More complicated 2-class problem $(4 \times 4 \text{ checker board})$.



▶ The support vectors in this example:



Solving the SVM optimization problem

- So far we have not considered any algorithms for solving for the SVM.
- We have to solve a constrained optimization problem to obtain the Lagrange multipliers and hence the SVM.
- Many specialized algorithms have been proposed for this.

▶ The optimization problem to be solved is

$$\max_{\pmb{\mu}} \qquad q(\pmb{\mu}) = \sum_{i=1}^n \, \mu_i - \frac{1}{2} \sum_{i,j=1}^n \, \mu_i \mu_j y_i y_j K(X_i, X_j)$$
 subject to
$$0 \leq \mu_i \leq C, \quad i=1,\dots,n, \quad \sum_i y_i \mu_i = 0$$

- ► A quadratic programming (QP) problem with interesting structure.
- ▶ Due to the special structure, many efficient algorithms are proposed.

- ▶ One interesting idea Chunking
- We optimize on only a few variables at a time.
- Dimensionality of the optimization problem is controlled.
- We keep randomly choosing the subset of variables.
- Gave rise to the first specialized algorithm for SVM SVM Light

SMO Algorithm

- ► Taking chunking to extreme level what is the smallest set of variables we can optimize on?
- We need to consider at least two variables because there is an equality constraint.
- Sequential Minimal Optimization (SMO) works on optimizing two variables at a time.
- We can analytically find the optimum with respect to two variables.
- ► The algorithm (heuristically) decides which two we consider in each iteration.
- A very efficient algorithm
- ► There are many such algorithms.

Kernel Trick

- ▶ We use $\phi: \Re^n \to \mathcal{H}$ to map pattern vectors into appropriate high dimensional space.
- ▶ Kernel fuction allows us to compute innerproducts in \mathcal{H} implicitly without using (or even knowing) ϕ .
- ► Through kernel functions, we learn nonlinear classifiers using 'linear techniques'.
- ▶ Algorithms that use only innerproducts (e.g., Fisher discriminant, regression etc) can be implicitly executed in a high dimensional, \mathcal{H} , by using a kernel function.
- We can elegantly construct non-linear versions of linear techniques.

Support Vector Regression

- Now we consider the regression problem.
- Given training data

$$\{(X_1, y_1), \dots, (X_n, y_n)\}, X_i \in \mathbb{R}^m, y_i \in \mathbb{R},$$

want to find 'best' function to predict y given X.

▶ We search in a parameterized class of functions

$$g(X, W) = w_1 \phi_1(X) + \dots + w_{m'} \phi_{m'}(X) + b$$

= $W^T \Phi(X) + b$,

where $\phi_i: \Re^m \to \Re$ are some chosen functions.

- ▶ If we choose, $\phi_i(X) = x_i$ (and hence, m = m') then it is the usual linear model.
- We are essentially learning a linear model in terms of $\Phi(X)$.
- ▶ We want to formulate the problem so that we can use the Kernel idea.
- Then, by using a kernel function, we never need to compute or even precisely specify the mapping Φ.

Loss function

lacktriangle In a general regression problem, we need to find W to minimize

$$\sum_{i} L(y_i, g(X_i, W))$$

where L is a loss function.

We consider a special loss function that allows us to use the kernel trick.

ϵ -insensitive loss

• We employ ϵ -insensitive loss function:

$$L_{\epsilon}(y_i, g(X_i, W)) = 0$$
 If $|y_i - g(X_i, W)| < \epsilon$
= $|y_i - g(X_i, W)| - \epsilon$ otherwise

Here, ϵ is a parameter of the loss function.

- ▶ If prediction is within ϵ of true value, there is no loss.
- Using absolute value of error rather than square of error allows for better robustness.
- Also gives us optimization problem with the right structure.
- ▶ Empirical risk minimization under the ϵ -insensitive loss function would minimize

$$\sum_{i=1}^{n} \max \left(|y_i - \Phi(X_i)^T W - b| - \epsilon, \ 0 \right)$$

 \blacktriangleright We want W, b to minimize

$$\sum_{i=1}^{n} \max \left(|y_i - \Phi(X_i)^T W - b| - \epsilon, \ 0 \right)$$

▶ We can pose the problem as follows.

$$\begin{aligned} & \underset{W,b,\boldsymbol{\xi},\boldsymbol{\xi'}}{\min} & & \sum_{i=1}^n \, \xi_i + \sum_{i=1}^n \, \xi_i' \\ & \text{subject to} & & y_i - W^T \Phi(X_i) - b \leq \epsilon + \xi_i, \quad i = 1, \dots, n \\ & & & W^T \Phi(X_i) + b - y_i \leq \epsilon + \xi_i', \quad i = 1, \dots, n \\ & & & \xi_i \geq 0, \ \xi_i' \geq 0 \quad i = 1, \dots, n \end{aligned}$$

- This does not give a dual with the structure we want.
- ▶ So, we reformulate the optimization problem

The Optimization Problem

▶ Find W, b and ξ_i, ξ'_i to

$$\begin{split} \text{minimize} & \quad \frac{1}{2}W^TW + C\left(\sum_{i=1}^n \ \xi_i + \sum_{i=1}^n \ \xi_i'\right) \\ \text{subject to} & \quad y_i - W^T\Phi(X_i) - b \leq \epsilon + \xi_i, \quad i = 1, \dots, n \\ & \quad W^T\Phi(X_i) + b - y_i \leq \epsilon + \xi_i', \quad i = 1, \dots, n \\ & \quad \xi_i \geq 0, \ \xi_i' \geq 0 \quad i = 1, \dots, n \end{split}$$

• We have added the term W^TW in the objective function. This is like model complexity in a regularization context.

- ▶ Like earlier, we can form the Lagrangian and then, using Kuhn-Tucker conditions, can get the optimal values of W and b.
- Given that this problem is similar to the earlier one, we would get W^* in terms of the optimal lagrange multipliers as earlier.
- Essentially, the lagrange multipliers corresponding to the inequality constraints on the errors would be the determining factors.
- ▶ We can use the same technique as earlier to formulate the dual to solve for the optimal Lagrange multipliers.

The dual

▶ The dual of this problem is

$$\begin{split} \max_{\pmb{\alpha},\pmb{\alpha}} \qquad & \sum_{i=1}^n \, y_i(\alpha_i - \alpha_i') - \epsilon \sum_{i=1}^n \, (\alpha_i + \alpha_i') \\ & - \frac{1}{2} \sum_{i,j} \, (\alpha_i - \alpha_i') (\alpha_j - \alpha_j') \Phi(X_i)^T \Phi(X_j) \\ \text{subject to} \qquad & \sum_{i=1}^n \, (\alpha_i - \alpha_i') = 0 \\ & 0 \leq \alpha_i, \, \, \alpha_i' \leq C, \quad i = 1, \dots, n \end{split}$$

▶ Here α_i and α'_i are the Lagrange multipliers corresponding to the first two inequalities in the primal.

The solution

- ▶ We can use the Kuhn-Tucker conditions to derive the final optimal values of *W* and *b* as earlier.
- This gives us

$$W^* = \sum_{i=1}^{n} (\alpha_i^* - \alpha_i^{*'}) \Phi(X_i)$$

$$b^* = y_j - \Phi(X_j)^T W^* + \epsilon, \quad j \text{ s.t. } 0 < \alpha_j^* < C/n$$

We have

$$W^* = \sum_{i=1}^{n} (\alpha_i^* - \alpha_i^{*'}) \Phi(X_i)$$

$$b^* = y_j - \Phi(X_j)^T W^* + \epsilon, \quad j \ s.t. \ 0 < \alpha_j^* < C/n$$

- Note that we have $\alpha_i^* \alpha_i^{*'} = 0$. Also, α_i^* , $\alpha_i^{*'}$ are zero for examples where error is less than ϵ .
- ► The final W is a linear combination of some of the examples – the support vectors.
- Note that the dual and the final solution are such that we can use the kernel trick.

- Let $K(X, X') = \Phi(X)^T \Phi(X')$.
- ▶ The optimal model learnt is

$$g(X, W^*) = \sum_{i=1}^{n} (\alpha_i^* - \alpha_i^{*'}) \phi(X_i)^T \phi(X) + b^*$$
$$= \sum_{i=1}^{n} (\alpha_i^* - \alpha_i^{*'}) K(X_i, X) + b^*$$

As earlier, b^* can also be written in terms of the Kernel function.

Support vector regression

- Once again, the kernel trick allows us to learn non-linear models using a linear method.
- ▶ The parameters: C, ϵ and parameters of kernel function.
- ► The basic idea of SVR can be used in many related problems.

SV regression

- ▶ With the ϵ -insensitive loss function, points whose targets are within ϵ of the prediction do not contribute any 'loss'.
- ▶ Gives rise to some interesting robustness of the method. It can be proved that local movements of target values of points outside the ϵ -tube do not influence the regression.
- ▶ Robustness essentially comes through the support vector representation of the regression.

- ▶ In our formulation of the regression problem we did not explain why we added W^TW term in the objective function.
- We are essentially minimizing

$$\frac{1}{2}W^{T}W + C \sum_{i=1}^{n} \max(|y_{i} - \Phi(X_{i})^{T}W - b| - \epsilon, 0)$$

- This is 'regularized risk minimization'.
- ▶ Then W^TW is the model complexity term which is intended to favour learning of 'smoother' models.
- Next we explain why W^TW is a good term to capture degree of smoothness in case of linear models.

- ▶ Let $f: \Re^m \to \Re$ be a continuous function.
- ▶ Continuity means we can make |f(X) f(X')| as small as we want by taking ||X X'|| sufficiently small.
- ► There are ways to characterize the 'degree of continuity' of a function.
- ▶ We consider one such measure now.

ϵ -Margin of a function

▶ The ϵ -margin of a function, $f: \Re^n \to \Re$ is

$$m_{\epsilon}(f) = \inf\{||X - X'|| : |f(X) - f(X')| \ge 2\epsilon\}$$

▶ The intuitive idea is:

How small can
$$||X - X'||$$
 be, still keeping $|f(X) - f(X')|$ 'large'

▶ The larger $m_{\epsilon}(f)$, the smoother is the function.

$$m_{\epsilon}(f) = \inf\{||X - X'|| : |f(X) - f(X')| \ge 2\epsilon\}$$

- ▶ Obviously, $m_{\epsilon}(f) = 0$ if f is discontinuous.
- $m_{\epsilon}(f)$ can be zero even for continuous functions, e.g., f(x)=1/x.
- $m_{\epsilon}(f) > 0$ for all $\epsilon > 0$ iff f is uniformly continuous.
- ► Higher margin would mean the function is 'slowly varying' and hence is a 'smoother' model.

Linear Models and margin

Consider regression with linear models. Then,

$$|f(X) - f(X')| = |W^T(X - X')|.$$

- ▶ For all X, X' with $|W^T(X X')| \ge 2\epsilon$, we want the smallest ||X X'||
- It would be smallest if $|W^T(X-X')|=2\epsilon \text{ and } (X-X') \text{ is parallel to } W.$ That is, $X-X'=\pm \frac{2\epsilon W}{WTW}.$
- ▶ Thus, $m_{\epsilon}(f) = || \pm \frac{2\epsilon W}{W^T W}|| = \frac{2\epsilon}{||W||}$.
- ▶ Thus in our optimization problem adding the term W^TW promotes learning of smoother models.
- ▶ As we have seen linear regression models use this as the regularization term.

- ► The basic idea of kernel functions, as we saw in SVM, has been extended in many ways.
- ► There have been many extensions of the basic SVM method also.
- ▶ Some of them are essentially formulations of approximate solutions to make the algorithm more efficient.
- Some of them are reformulations to add additional features to the SVM method.
- We consider a couple of simple examples of such extensions.

▶ Suppose the optimization problem is changed to

$$\begin{aligned} & \underset{W,b,\pmb{\xi}}{\min} & & \frac{1}{2}W^TW \,+\, b^2 \,+\, C\sum_{i=1}^n \,\xi_i \\ & \text{subject to} & & y_i(W^TX_i+b) \geq 1-\xi_i, \quad i=1,\dots,n \\ & & \xi_i \geq 0, \quad i=1,\dots,n \end{aligned}$$

We have added the b^2 term to the objective function. The main reason is that it simplifies the dual. ▶ The dual turns out to be

$$\begin{aligned} \max_{\pmb{\mu}} \qquad \sum_{i=1}^n \; \mu_i \; - \; \frac{1}{2} \sum_{i,j=1}^n \; \mu_i \mu_j y_i y_j K(X_i, X_j) \\ - \; \frac{1}{2} \sum_{i,j=1}^n \; \mu_i \mu_j y_i y_j \\ \text{subject to} \qquad 0 \leq \mu_i \leq C, \quad i=1,\dots,n, \end{aligned}$$

- The equality constraint is absent.
 Only bound constraints on variables.
- Allows for efficient optimization. (Successive overrelaxation).

- Next, we consider a reformulation of SVM optimization problem, known as ν -SVM.
- Recall that the primal problem for SVM with slack variables is

$$\begin{aligned} & \min_{W,b,\pmb{\xi}} & & \frac{1}{2}W^TW \,+\, C\sum_{i=1}^n\,\xi_i \\ & \text{subject to} & & y_i(W^T\phi(X_i)+b) \geq 1-\xi_i, \ i=1,\dots,n \\ & & \xi_i \geq 0, \ i=1,\dots,n \end{aligned}$$

ν -SVM

- ▶ In the SVM formulation with slack variables, we do not know how to choose *C*.
- Consider a changed optimization problem

$$\begin{aligned} \min_{W,b,\pmb{\xi},\rho} & & \frac{1}{2}W^TW - \nu\rho + \frac{1}{n}\sum \xi_i \\ \text{subject to} & & y_i[W^T\phi(X_i) + b] \geq \rho - \xi_i \\ & & \xi_i \geq 0. \end{aligned}$$

where ν is a user-chosen constant.

- ▶ Note that $W, b, \rho, \xi_i = 0$ is a feasible solution.
- ▶ We do not need $\rho \ge 0$ constraint.

▶ The Lagrangian for this problem is

$$L(W, b, \xi, \rho, \eta, \mu) = \frac{1}{2} W^T W - \nu \rho + \frac{1}{n} \sum_{i=1}^n \xi_i$$
$$- \sum_{i=1}^n \eta_i \xi_i + \sum_{i=1}^n \mu_i \left(\rho - \xi_i - y_i [W^T \phi(X_i) + b] \right)$$

▶ The μ_i are the Lagrange multipliers for the separability constraints and η_i are the Lagrange multipliers for the constraints $\xi_i \geq 0$.

The Kuhn-Tucker conditions give us

$$\blacktriangleright \nabla_W L = 0 \Rightarrow W = \sum_i \mu_i y_i \phi(X_i)$$

$$ightharpoonup \frac{\partial L}{\partial b} = 0 \Rightarrow \sum \mu_i y_i = 0$$

$$\blacktriangleright \frac{\partial L}{\partial a} = 0 \Rightarrow \sum \mu_i = \nu$$

$$\rho - \xi_i - y_i(W^T \phi(X_i) + b) \le 0; \quad \xi_i \ge 0; \quad \forall i$$

$$\mu_i \ge 0; \ \eta_i \ge 0, \ \forall i$$

$$\mu_i(\rho - \xi_i - y_i(W^T\phi(X_i) + b)) = 0; \quad \eta_i \xi_i = 0, \ \forall i$$

▶ Suppose $\xi_i > 0$ for some i. Then we have $\eta_i = 0$ and hence $\mu_i = \frac{1}{n}$. Hence

$$\nu = \sum_{i=1}^{n} \mu_{i} = \sum_{i:\xi_{i}>0} \mu_{i} + \sum_{i:\xi_{i}=0} \mu_{i}$$

$$\geq \sum_{i:\xi_{i}>0} \mu_{i} = \frac{|\{i:\xi_{i}>0\}|}{n}$$

Hence we have:
 ν is an upper bound on the fraction of 'margin errors'.

• We also have, because $0 \le \mu_i \le \frac{1}{n}$,

$$\nu = \sum_{i=1}^{n} \mu_{i} = \sum_{i:\mu_{i}>0} \mu_{i} + \sum_{i:\mu_{i}=0} \mu_{i}$$

$$\leq \sum_{i:\mu_{i}>0} \mu_{i} \leq \frac{|\{i:\mu_{i}>0\}|}{n}$$

Hence we have:
 ν is a lower bound on the fraction of support vectors.

- ▶ In the ν -SVM formulation, the ν is the user chosen constant.
- Unlike the parameter C, the ν has an interesting interpretation.
- ▶ It is simultaneously the upperbound on fraction of errors and lower bound on fraction of support vectors.
- ▶ If for the chosen ν , the problem has a solution with $\rho > 0$, then the bounds would be met.
- ► This gives us a good way to choose this 'penalty constant'.

▶ The dual for the ν -SVM turns out to be

$$\max_{\pmb{\mu}} \qquad q(\pmb{\mu}) = -\frac{1}{2} \sum_{i,j=1}^n \ \mu_i \mu_j y_i y_j K(X_i, X_j)$$
 subject to
$$0 \leq \mu_i \leq \frac{1}{n}, \forall i; \ \sum_{i=1}^n \ y_i \mu_i = 0; \ \sum_{i=1}^n \mu_i = \nu$$

- This a simple optimization problem similar to that of 'C-SVM'.
- ▶ One can show that if we have a solution for ν -SVM then if we choose $C=1/\rho n$, we get the same solution with 'C-SVM'.

ν SVR

- ▶ This idea can be extended to the regression problem also.
- ▶ In support vector regression, we had two user defined constants: ϵ and C.
- ▶ The ϵ specifies the 'tolerable error' and it is difficult to know what value to choose for it.
- We can reformulate SVR so that we can optimize on ϵ also.
- ▶ This will be very similar to the ν -SVM formulation.

▶ Recall the optimization problem in SVR:

$$\begin{aligned} & \underset{W,b,\boldsymbol{\xi},\boldsymbol{\xi'}}{\min} & & \frac{1}{2}W^TW + C\left(\sum_{i=1}^n \ \xi_i + \sum_{i=1}^n \ \xi_i'\right) \\ & \text{subject to} & & y_i - W^T\Phi(X_i) - b \leq \epsilon + \xi_i, \quad i = 1,\dots,n \\ & & & W^T\Phi(X_i) + b - y_i \leq \epsilon + \xi_i', \quad i = 1,\dots,n \\ & & & \xi_i \geq 0, \ \xi_i' \geq 0 \quad i = 1,\dots,n \end{aligned}$$

▶ We change the optimization problem to the following:

$$\begin{aligned} & \underset{W,b,\epsilon,\boldsymbol{\xi},\boldsymbol{\xi'}}{\min} & & \frac{1}{2}W^TW + C\left(\nu\epsilon + \frac{1}{n}\sum_{i=1}^n\left(\xi_i + \xi_i'\right)\right) \\ & \text{subject to} & & y_i - W^T\phi(X_i) - b \leq \epsilon + \xi_i, \quad i = 1,\dots,n \\ & & & W^T\phi(X_i) + b - y_i \leq \epsilon + \xi_i', \quad i = 1,\dots,n \\ & & & \xi_i \geq 0, \; \xi_i' \geq 0 \; \epsilon \geq 0, \; i = 1,\dots,n \end{aligned}$$

where ν is a user-chosen constant.

• We get similar results as in ν -SVM.

Risk minimization view of SVM

- We posed the support vector regression problem as a (regularized) risk minimization under a special loss function.
- It was then reformulated into an (equivalent) constrained optimization problem.
- ▶ In contrast, we formulated the SVM directly as a constrained optimization problem.
- However, it can also be seen to be minimization of (regularized) empirical risk under a special loss function.

▶ The optimization problem for SVM is

$$\begin{aligned} & \underset{W,b,\pmb{\xi}}{\min} & & \frac{1}{2}W^TW \,+\, C\sum_{i=1}^n\,\xi_i \\ & \text{subject to} & & y_i(W^TX_i+b) \geq 1-\xi_i, \quad i=1,\dots,n \\ & & \xi_i \geq 0, \quad i=1,\dots,n \end{aligned}$$

• Given any W, b, the ξ_i have to satisfy

$$\xi_i \ge \max(0, 1 - y_i(W^T X_i + b))$$

▶ Since we need to minimize $\sum \xi_i$, we need to take the value above for each ξ_i .

► Hence we can find SVM by solving the following unconstrained optimization problem:

$$\min_{W,b} \ \frac{1}{2} W^T W + C \sum_{i=1}^n \max(0, 1 - y_i (W^T X_i + b))$$

- The model (or classifier) we are learning is $f(X) = W^T X + b$.
- For this model, we already saw W^TW is a good regularization term.

Consider the loss function defined by

$$L_{\mathsf{hinge}}(y, f(X)) = \max(0, 1 - yf(X))$$

▶ Then the optimization problem is same as

$$\min_{W,b} \frac{1}{n} \sum_{i=1}^{n} L(y_i, f(X_i)) + C' \frac{1}{2} W^T W$$

► Thus, our SVM formulation is empirical risk minimization under hinge-loss along with a regularization term.

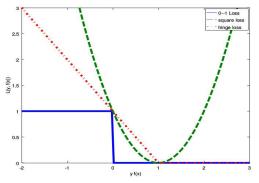
- ► As we saw earlier, the hinge-loss and square-loss are good convex approximations of the 0–1 loss.
- ▶ For 0–1 loss L(y, f(X)) is one if yf(X) is negative and zero otherwise.
- ▶ The squared error loss can be written as

$$L_{\mathsf{square}}(y, h(X)) = (1 - yf(X))^2$$

▶ The hinge loss is given by

$$L_{\mathsf{hinge}}(y, h(X)) = \max(0, \ 1 - yf(X))$$

▶ We can plot all the functions as follows.



(Here we plot yf(X) on x-axis and L(y,f(X)) on y-axis).

- Hinge loss is also called soft-margin loss.
- ightharpoonup Supoose we want to minimize, over all f,

$$E[\max(0, 1 - yf(X))], y \in \{+1, -1\}$$

- Intuitively the best we can do is to make sign of f(X) to be same as sign of the corresponding y.
- \blacktriangleright Hence, intuitively, the best f is

$$f(X) > 0$$
, if $P[y = +1|X] > 0.5$; else $f(X) < 0$

This is indeed a good classifier.

- ▶ In SVM method, there are two important ingradients.
- One is the Kernel function.
- Kernel functions allow us to learn nonlinear models using essentially linear techniques.
- Second is the 'support vector' expansion the final model is expressed as a ('sparse') linear combination of some of the data vectors.
- Kernels are a good way to capture 'similarity' and are useful in general.
- ► The support vector expansion is also a general property of Kernel based methods.
- ▶ We look at this general view of Kernels next.