Support vector machines

Pavlo Mozharovskyi¹

¹LTCI, Télécom Paris, Institut Polytechnique de Paris

Machine learning

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Today

Vapnik-Chervonenkis theory, simplest case

The support vector machine

Optimal margin classifier Introducing kernels (1992)

Allowing for misclassification: soft margin (1995)

Implementation

Literature

Learning materials include but are not limited to:

- Hastie, T., Tibshirani, R., and Friedman, J. (2009). The Elements of Statistics Learning: Data Mining, Inference, and Prediction (Second Edition). Springer.
 - ► Section 12.{1,2,3.1,3.2}.
- Slides of the lecture.
- Boser, B. E., Guyon, I. M., and V. N. Vapnik (1992).
 A training algorithm for optimal margin classifiers.
 In: Proceedings of the Fifth Annual Workshop of Computational Learning Theory, Pittsburgh, ACM, 5, 144–152.
- Cortes, C. and Vapnik, V. (1995).
 Support-vector networks.
 Machine learning, 20, 273–297.
- Vapnik, V. N. (1998).
 Statistical Learning Theory.
 John Wiley & Sons.



Binary supervised classification (reminder)

Notation:

- ▶ **Given:** for the random pair (X, Y) in $\mathbb{R}^d \times \{-1, 1\}$ consisting of a random observation X and its random binary label Y (class), a sample of n i.i.d.: $(\mathbf{x}_1, y_1), ..., (\mathbf{x}_n, y_n)$.
- **Goal:** predict the label of the new (unseen before) observation x.
- ▶ **Method:** construct a classification rule:

$$g: \mathbb{R}^d \to \{-1,1\}, \mathbf{x} \mapsto g(\mathbf{x}),$$

so g(x) is the prediction of the label for observation x.

► **Criterion:** of the performance of *g* is the **error probability**:

$$R(g) = \mathbb{P}[g(X) \neq Y] = \mathbb{E}[\mathbb{1}(g(X) \neq Y)].$$

▶ **The best solution:** is to know the distribution of (X,Y):

$$g(\mathbf{x}) = \operatorname{sign}(2\mathbb{E}[Y|X=\mathbf{x}] - 1 > 0)$$
.



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What can be said about the **error probability**, *i.e.* about the relationship between

$$\mathbb{P}\big(g(X, \boldsymbol{w}) \neq Y\big) = \int_{\mathbb{R}^d} \mathbb{1}\big(g(\boldsymbol{x}, \boldsymbol{w}) \neq Y\big) dF_X \quad \text{and} \quad \frac{1}{n} \sum_{i=1}^n \mathbb{1}\big(g(\boldsymbol{x}_i, \boldsymbol{w}) \neq y_i\big) ?$$

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Let $X_1, ..., X_n$ be a random sample on \mathbb{R} . The **empirical distribution** function is defined as

$$\mathbb{F}_n(t) = \frac{1}{n} \sum \mathbb{1}(X_i \leq t).$$

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Theorem (Glivenko-Cantelli)

If $X_1, X_2, ...$ are i.i.d. random variables with distribution function F, then

$$\|\mathbb{F}_n - F\|_{\infty} = \sup_{x \in \mathbb{R}} |\mathbb{F}_n(x) - F(x)| \xrightarrow{a.s.} 0.$$

Under additional conditions, for $g(\mathbf{x}, \mathbf{w})$ and a probability measure F_X , for any $\epsilon > 0$ it holds

$$\mathbb{P}\Big\{\sup_{\boldsymbol{w}}\Big(\underbrace{\mathbb{P}\big(g(X,\boldsymbol{w})\neq Y\big)}_{L\big(g(\cdot,\boldsymbol{w})\big)}-\underbrace{\frac{1}{n}\sum_{i=1}^{n}\mathbb{1}\big(g(X_{i},\boldsymbol{w})\neq Y_{i}\big)}_{L_{emp}\big(g(\cdot,\boldsymbol{w})\big)}\Big)>\epsilon\Big\}\underset{n\to\infty}{\longrightarrow}0.$$

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$$\mathbb{P}\left\{\sup_{k\in\{1,\ldots,N\}}\left(L(g(\cdot,\boldsymbol{w}_{k}))-L_{emp}(g(\cdot,\boldsymbol{w}_{k}))\right)>\epsilon\right\}$$

$$\leq \sum_{k=1}^{N}\mathbb{P}\left\{\left(L(g(\cdot,\boldsymbol{w}_{k}))-L_{emp}(g(\cdot,\boldsymbol{w}_{k}))\right)>\epsilon\right\}$$

Theorem (Chernoff-Hoeffding, Bernoulli scheme)

If $X_1,...,X_n$ are i.i.d. random variables taking values in $\{0,1\}$, then for any $\epsilon>0$ it holds

$$\mathbb{P}\Big(\mathbb{E}[X_i] - \frac{1}{n}\sum_{i=1}^n X_i > \epsilon\Big) < e^{-2\epsilon^2 n}.$$

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$$\begin{split} & \sum_{k=1}^{N} \mathbb{P} \Big\{ \Big(L \big(g(\cdot, \boldsymbol{w}_{k}) \big) - L_{emp} \big(g(\cdot, \boldsymbol{w}_{k}) \big) \Big) > \epsilon \Big\} \\ &= \sum_{k=1}^{N} \mathbb{P} \Big\{ \Big(\underbrace{\mathbb{P} \big(g(X, \boldsymbol{w}_{k}) \neq Y \big)}_{\mathbb{E} \big[\mathbb{1} \big(g(X, \boldsymbol{w}_{k}) \neq Y \big) \big]} - \frac{1}{n} \sum_{i=1}^{n} \mathbb{1} \big(g(X_{i}, \boldsymbol{w}_{k}) \neq Y_{i} \big) \Big) > \epsilon \Big\} \end{split}$$

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$$= \sum_{k=1}^{N} \mathbb{P}\left\{ \left(\underbrace{\mathbb{P}(g(X, \boldsymbol{w}_{k}) \neq Y)}_{\mathbb{E}\left[\mathbb{1}\left(g(X, \boldsymbol{w}_{k}) \neq Y\right)\right]} - \frac{1}{n} \sum_{i=1}^{n} \mathbb{1}\left(g(X_{i}, \boldsymbol{w}_{k}) \neq Y_{i}\right) \right) > \epsilon \right\}$$

$$\leq Ne^{-2\epsilon^{2}n}.$$

So:

$$\mathbb{P}\Big\{\sup_{k\in\{1,\dots,N\}}\Big(L\big(g(\cdot,\boldsymbol{w}_k)\big)-L_{emp}\big(g(\cdot,\boldsymbol{w}_k)\big)\Big)>\epsilon\Big\}\leq Ne^{-2\epsilon^2n}\,.$$

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Let us fix this probability having chosen 0 < $\eta \leq$ 1, by that maintaining reliability 1 - η :

$$Ne^{-2\epsilon^2 n} = \eta$$
 or equivalently $\epsilon = \sqrt{\frac{\log N - \log \eta}{2n}}$.

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Theorem (Vapnik-Chervonenkis, 1974)

If from a set consisting of N classification rules a rule $g(\cdot, \mathbf{w})$ is chosen, which delivers empirical risk $L_{emp}(g(\cdot, \mathbf{w}))$, then with reliability $1-\eta$ one can state that the error probability $L(g(\cdot, \mathbf{w}))$ is bounded from above as follows

$$L(g(\cdot, \boldsymbol{w})) \leq L_{emp}(g(\cdot, \boldsymbol{w})) + \sqrt{\frac{\log N - \log \eta}{2n}}$$
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The number $\Phi(d, n)$ of all possible separations of n points in \mathbb{R}^d by a hyperplane via the origin is computed as

$$\Phi(d,n) = \begin{cases} 2\sum_{l=0}^{d-1} \binom{n-1}{l} & \text{if } d \leq n, \\ 2^n & \text{otherwise}. \end{cases}$$

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Plugging this into the Vapnik-Chervonenkis inequality gives:

$$L(g(\cdot, \boldsymbol{w})) \leq L_{emp}(g(\cdot, \boldsymbol{w})) + \sqrt{\frac{d \log n - \log \eta}{2n}}$$
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- As an example, consider the case of two linearly separable training classes. Even in this case, only little can be said about probability of points from one class inside the other one.
- Sticking to this "trivial" case, the safest separating hyperplane would be the one having maximal and equal margin to each of the classes.
- ► Finding such a hyperplane in a systematic way constitutes the main idea of the **optimal margin hyperplane** algorithm.

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Optimal margin hyperplane

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- ▶ This set is said to be **linearly separable** if there exist a non-zero vector $\psi \in \mathbb{R}^d$ and a scalar $b \in \mathbb{R}$ such that the n following inequalities hold:

$$\psi^T \mathbf{x}_i + b \ge 0$$
 if $y_i = 1$,
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- ▶ Involving the output (in this notation corresponding to the sign) allows for rewriting the above (restricting) inequalities in the following way:

$$\frac{y_i(\boldsymbol{\psi}^T\boldsymbol{x}_i+b)}{\|\boldsymbol{\psi}\|} \geq M, \quad i=1,...,n.$$

ightharpoonup The objective of the training algorithm is then to find the parameter vector ψ that maximizes M:

$$M^* = \max M$$
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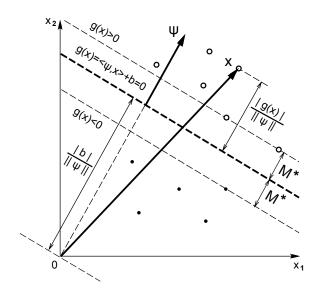
$$\min_{i \in \{1,\ldots,n\}} y_i(\boldsymbol{\psi}^T \boldsymbol{x}_i + b) = M^*.$$

- These patterns are called the support vectors of the decision boundary.
- ► Thus, the problem of finding a hyperplane with maximum margin can be seen as a **minimax** problem:

$$\max_{\boldsymbol{\psi} \in \mathbb{R}^d, \|\boldsymbol{\psi}\| = 1} \min_{i \in \{1, \dots, n\}} y_i(\boldsymbol{\psi}^T \boldsymbol{x}_i + b).$$



Optimal margin hyperplane: illustration



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- ▶ Then the problem of finding a maximum margin separating hyperplane, characterized by ψ , reduces to solving the following **quadratic optimization problem**:

$$\min \frac{1}{2} \| \boldsymbol{\psi} \|^2$$
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- This approach is impractical:
 - if the **dimension** d is large or **infinite**,
 - because no information about **support vectors** is gained.



► Construct a Lagrangian:

$$L(\boldsymbol{\psi}, b, \boldsymbol{\Lambda}) = \frac{1}{2} \boldsymbol{\psi}^T \boldsymbol{\psi} - \sum_{i=1}^n \alpha_i (y_i (\boldsymbol{\psi}^T \boldsymbol{x}_i + b) - 1)$$

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- ▶ The solution to the optimization problem is determined by the saddle point of this Lagrangian in the (d+1+n)-dimensional space of ψ , b, and Λ .
- ▶ The **minimum** should be taken w.r.t. the parameters ψ and b, the **maximum** should be taken w.r.t. the Lagrange multipliers Λ .

 \blacktriangleright At the point of minimum (w.r.t. ψ and b) one obtains:

$$\frac{\partial L(\psi, b, \mathbf{\Lambda})}{\partial \psi}\Big|_{\psi=\psi^*} = \left(\psi^* - \sum_{i=1}^n \alpha_i y_i \mathbf{x}_i\right) = 0,$$

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- ► This means that the optimal hyperplane can be written as a **linear** combination of training observations.
- ▶ Only training observations x_i with (strictly) positive Lagrange multipliers (*i.e.* with $\alpha_i > 0$) have an efficient contribution to the sum the **support vectors**.

Substitution of the minimum conditions into the Lagrangian yields the following optimization problem:

$$\max W(\mathbf{\Lambda}) = \sum_{i=1}^{n} \alpha_i - \frac{1}{2} \sum_{i=1}^{n} \sum_{j=1}^{n} \alpha_i \alpha_j y_i y_j \mathbf{x}_i^T \mathbf{x}_j$$

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Usually it is written in the matrix form:

$$\max W(\mathbf{\Lambda}) = \mathbf{\Lambda}^T \mathbf{1} - \frac{1}{2} \mathbf{\Lambda}^T \mathbf{D} \mathbf{\Lambda}$$

subject to
$$\mathbf{\Lambda}^T \mathbf{Y} = 0,$$

$$\mathbf{\Lambda} > \mathbf{0}$$

with \boldsymbol{D} being a $(n \times n)$ -dimensional matrix with entries $D_{ij} = y_i y_j \boldsymbol{x}_i^T \boldsymbol{x}_j$, $\boldsymbol{Y} = (y_1, ..., y_n)^T$, and $\boldsymbol{0}$ and $\boldsymbol{1}$ standing for n-dimensional vectors of zeros and ones.

Optimal margin hyperplane: classification

▶ After the optimal pair (ψ^*, b^*) is obtained, classification of an observation $\mathbf{x} \in \mathbb{R}^d$ reduces to determining its position in the projection onto ψ^* :

$$g(\mathbf{x}) = \operatorname{sign}(\psi^{*T} \mathbf{x} + b^{*})$$
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From this it becomes clear how to calculate b^* : it should position the separating hyperplane exactly in the middle between two support vectors from different classes, in the projection onto ψ^* :

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➤ Only **support vectors** influence the classification rule.

(Analogy with a mine field on the front line between two enemies.)



Optimal margin classifier (algorithm)

Finding the optimal margin hyperplane (training)

Input: Training sample $((x_1, y_1), ..., (x_n, y_n)) \subset \mathbb{R}^d \times \{-1, 1\}.$

1. Solve the constraint quadratic optimization problem to obtain $\mathbf{\Lambda}^* = (\alpha_1^*, ..., \alpha_n^*)^T$:

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Output: The classifier: $g(x) = \text{sign}\left(\sum_{i=1}^{n} y_i \alpha_i^* x_i^T x + b^*\right)$.



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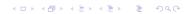
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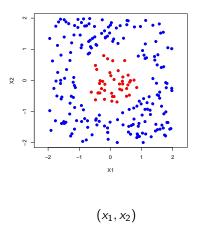
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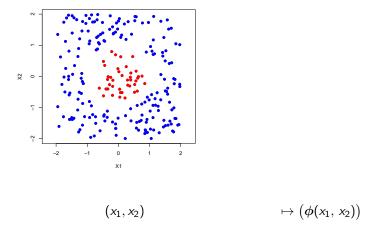
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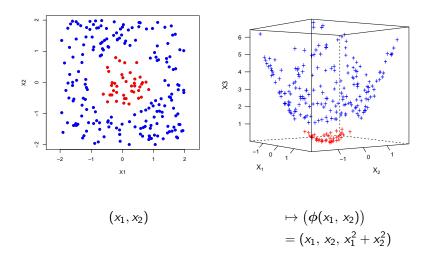
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Note: \mathbb{R}^D can be of infinite dimension.









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The quadratic problem depends only inner products as well.



▶ Consider the general form of the inner product in a Hilbert space:

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▶ A sufficient condition to ensure that K(u, v) defines an inner product in the feature space is that all the eigenvalues λ_i are positive.



▶ According to Mercer's theorem, for λ_i s to be positive, it is necessary and sufficient that

$$\int \int K(\boldsymbol{u},\boldsymbol{v})h(\boldsymbol{u})h(\boldsymbol{v})d\boldsymbol{u}d\boldsymbol{v}>0$$

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Polynomial kernel:

$$K(\boldsymbol{u}, \boldsymbol{v}) = (\boldsymbol{u}^T \boldsymbol{v} + 1)^{\beta}.$$



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▶ The **decision rule** has then form:

$$g(\mathbf{x}) = \sum_{i=1}^n y_i \alpha_i^* K(\mathbf{x}, \mathbf{x}_i) + b^*.$$

where one can only restrict to support vectors \mathbf{x}_i and their coefficients $\alpha_i^* > 0$.



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For sufficiently small σ the minimized functional describes the number of errors on the training set.



▶ In the minimum, strictly positive $\xi_{i_j} > 0$, j = 1, ..., k will identify the minimal subset of training errors:

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For sufficiently large C and sufficiently small σ , the pair (ψ^*, b^*) minimizing this functional will determine the **hyperplane** minimizing the number of errors and separating the rest with maximum margin.

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- ► For the reasons of computational tractability, we consider the (most commonly used) case:

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▶ The problem then becomes:

min
$$\frac{1}{2} \|\psi\|^2 + C \sum_{i=1}^n \xi_i$$
 subject to $y_i(\psi^T \mathbf{x}_i + b) \ge 1 - \xi_i$, $i = 1, ..., n$, $\xi_i \ge 0$, $i = 1, ..., n$.

▶ The corresponding Lagrangian is:

$$L(\boldsymbol{\psi}, b, \boldsymbol{\xi}, \boldsymbol{\Lambda}, \boldsymbol{r}) = \frac{1}{2} \boldsymbol{\psi}^T \boldsymbol{\psi} + C \sum_{i=1}^n \xi_i - \sum_{i=1}^n \alpha_i (y_i (\boldsymbol{\psi}^T \boldsymbol{x}_i + b) - 1 + \xi_i) - \sum_{i=1}^n r_i \xi_i$$

with $\mathbf{\Lambda} = (\alpha_1, ..., \alpha_n)^T$ and $\mathbf{r} = (r_1, ..., r_n)^T$ being the vectors of non-negative **Lagrange multipliers** corresponding to the two groups of inequality constraints.

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- The solution to the optimization problem is determined by the saddle point of this Lagrangian in the (d+1+n+n+n)-dimensional space of ψ , b, ξ , Λ , and r.
- ▶ The **minimum** should be taken w.r.t. the parameters ψ , b, and ξ , the **maximum** should be taken w.r.t. the Lagrange multipliers Λ and r.

▶ At the point of minimum (w.r.t. ψ , b, and ξ) one obtains:

$$\frac{\partial L(\psi, b, \xi, \mathbf{\Lambda}, \mathbf{r})}{\partial \psi}\Big|_{\psi=\psi^*} = \left(\psi^* - \sum_{i=1}^n \alpha_i y_i \mathbf{x}_i\right) = 0,$$

$$\frac{\partial L(\psi, b, \xi, \mathbf{\Lambda}, \mathbf{r})}{\partial b}\Big|_{b=b^*} = \sum_{i=1}^n y_i \alpha_i = 0,$$

$$\frac{\partial L(\psi, b, \xi, \mathbf{\Lambda}, \mathbf{r})}{\partial \xi_i}\Big|_{\xi_i = \xi_i^*} = C - \alpha_i - r_i = 0, \quad i = 1, ..., n.$$

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$$\begin{split} \frac{\partial L(\boldsymbol{\psi}, \boldsymbol{b}, \boldsymbol{\xi}, \boldsymbol{\Lambda}, \boldsymbol{r})}{\partial \boldsymbol{\psi}} \Big|_{\boldsymbol{\psi} = \boldsymbol{\psi}^*} &= \left(\boldsymbol{\psi}^* - \sum_{i=1}^n \alpha_i y_i \boldsymbol{x}_i \right) = 0, \\ \frac{\partial L(\boldsymbol{\psi}, \boldsymbol{b}, \boldsymbol{\xi}, \boldsymbol{\Lambda}, \boldsymbol{r})}{\partial \boldsymbol{b}} \Big|_{\boldsymbol{b} = \boldsymbol{b}^*} &= \sum_{i=1}^n y_i \alpha_i = 0, \\ \frac{\partial L(\boldsymbol{\psi}, \boldsymbol{b}, \boldsymbol{\xi}, \boldsymbol{\Lambda}, \boldsymbol{r})}{\partial \boldsymbol{\xi}_i} \Big|_{\boldsymbol{\xi}_i = \boldsymbol{\xi}_i^*} &= C - \alpha_i - r_i = 0, \quad i = 1, ..., n. \end{split}$$

▶ This leads to the following quadratic problem:

$$\begin{aligned} \max W(\mathbf{\Lambda}) &= \sum_{i=1}^n \alpha_i - \frac{1}{2} \sum_{i=1}^n \sum_{j=1}^n \alpha_i \alpha_j y_i y_j \mathbf{x}_i^T \mathbf{x}_j \\ \text{subject to} & \sum_{i=1}^n \alpha_i y_i = 0 \,, \\ & 0 \leq \alpha_i \leq C \,, \quad i = 1, ..., n \,. \end{aligned}$$

Support vector machine (SVM)

► The **training** phase:

max
$$\mathbf{\Lambda}^T \mathbf{1} - \frac{1}{2} \mathbf{\Lambda}^T \mathbf{D} \mathbf{\Lambda}$$
 subject to $\mathbf{\Lambda}^T \mathbf{Y} = 0$, $\mathbf{0} < \mathbf{\Lambda} < C\mathbf{1}$.

with $\mathbf{Y} = (y_1, ..., y_n)^T$, $\mathbf{0}$ and $\mathbf{1}$ standing for *n*-dimensional vectors of zeros and ones, C being a properly chosen constant, and \mathbf{D} being a $(n \times n)$ -dimensional matrix with entries

$$D_{ij} = y_i y_j K(\mathbf{x}_i, \mathbf{x}_j), \quad i, j = 1, ..., n,$$

where $K(\boldsymbol{u}, \boldsymbol{v})$ is a properly chosen kernel function. The result is the optimal vector $\boldsymbol{\Lambda}^* = (\alpha_1^*, ..., \alpha_n^*)^T$.

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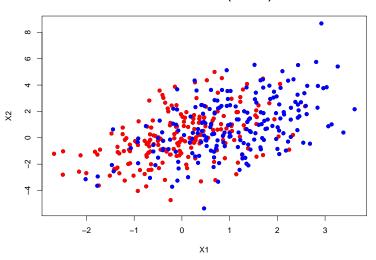
where $K(\boldsymbol{u}, \boldsymbol{v})$ is a properly chosen kernel function. The result is the optimal vector $\boldsymbol{\Lambda}^* = (\alpha_1^*, ..., \alpha_n^*)^T$.

Then, taking any two support vectors \mathbf{x}_{i_A} and \mathbf{x}_{i_B} from opposite classes, *i.e.* with $i_A \in \arg\max_{j: y_j = 1, \alpha_j^* > 0} \sum_{i=1}^n y_i \alpha_i^* K(\mathbf{x}_j, \mathbf{x}_i)$ and $i_B \in \arg\min_{j: y_j = -1, \alpha_i^* > 0} \sum_{i=1}^n y_i \alpha_i^* K(\mathbf{x}_j, \mathbf{x}_i)$, calculate threshold:

$$b^* = -\frac{1}{2} \sum_{i=1}^n y_i \alpha_i^* \left(K(\mathbf{x}_i, \mathbf{x}_{i_A}) + K(\mathbf{x}_i, \mathbf{x}_{i_B}) \right).$$

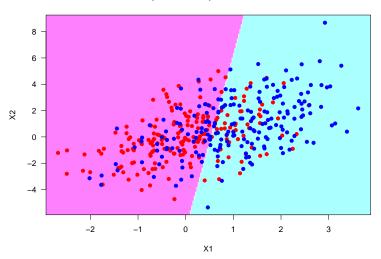
Normal location alternative





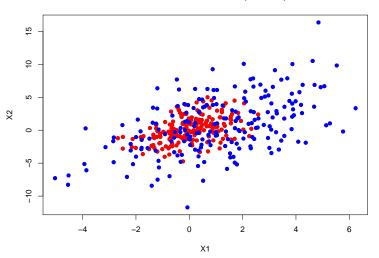
SVM: normal location alternative

SVM (linear kernel) for Normal1 data



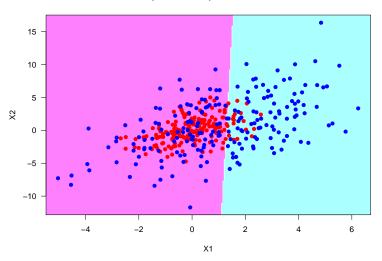
Normal location-scale alternative



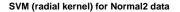


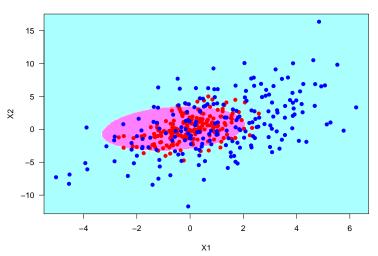
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For R-software, SVM is implemented in such packages as, e.g., e1071, kernlab, klaR, sympath.

For an overview, see, e.g.:

Karatzoglou, A., Meyer, D., and Hornik, K. (2006).
 Support vector machines in R.
 Journal of Statistical Software, 15(9).

Thank you for your attention!

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And some references

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