Machine Learning and Bayesian Inference

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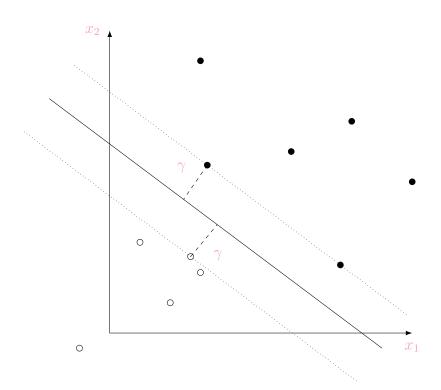
Part II

Support vector machines

General methodology

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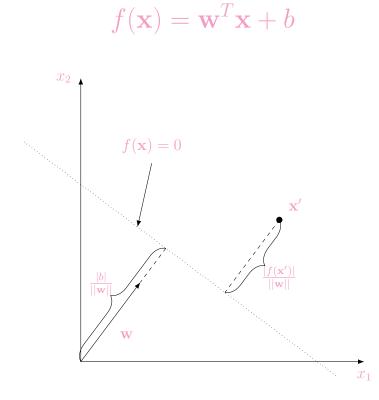
Suggestion: why not drop all this probability nonsense and just do this:



Draw the boundary as far away from the examples as possible.

The distance γ is the *margin*, and this is the *maximum margin classifier*.

If you completed the *exercises for AII* then you'll know that linear classifiers have a very simple geometry. For



For x' on one side of the line $f(\mathbf{x}) = 0$ we have $f(\mathbf{x}') > 0$ and on the other side $f(\mathbf{x}') < 0$.

Problems:

- Given the usual training data s, can we now find a *training algorithm* for obtaining the weights?
- What happens when the data is not *linearly separable*?

To derive the necessary training algorithm we need to know something about *constrained optimization*.

We can address the second issue with a simple modification. This leads to the Support Vector Machine (SVM).

Despite being decidedly "non-Bayesian" the SVM is currently a gold-standard:

Do we need hundreds of classifiers to solve real world classification problems, Fernández-Delgardo at al., Journal of Machine Learning Research 2014.

You are familiar with *maximizing* and *minimizing* a function $f(\mathbf{x})$. This is *unconstrained optimization*.

We want to extend this:

- 1. Minimize a function $f(\mathbf{x})$ with the constraint that $g(\mathbf{x}) = 0$.
- 2. Minimize a function $f(\mathbf{x})$ with the constraints that $g(\mathbf{x}) = 0$ and $h(\mathbf{x}) \ge 0$.

Ultimately we will need to be able to solve problems of the form: find \mathbf{x}_{opt} such that

 $\mathbf{x}_{opt} = \operatorname*{argmin}_{\mathbf{x}} f(\mathbf{x})$

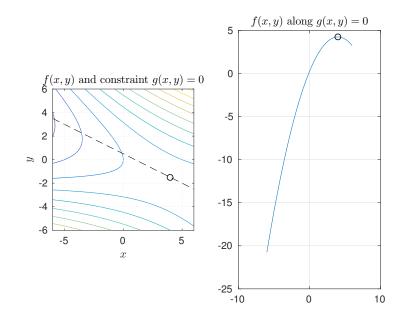
under the constraints

 $g_i(\mathbf{x}) = 0$ for i = 1, 2, ..., n

and

 $h_j(\mathbf{x}) \ge 0$ for j = 1, 2, ..., m.

For example:



Minimize the function

$$f(x,y) = -\left(2x + y^2 + xy\right)$$

subject to the constraint

$$g(x, y) = x + 2y - 1 = 0.$$

Step 1: introduce the *Lagrange multiplier* λ and form the *Langrangian* $L(x, y, \lambda) = f(x, y) - \lambda g(x, y)$

Necessary condition: it can be shown that if (x', y') is a solution then $\exists \lambda'$ such that



So for our example we need

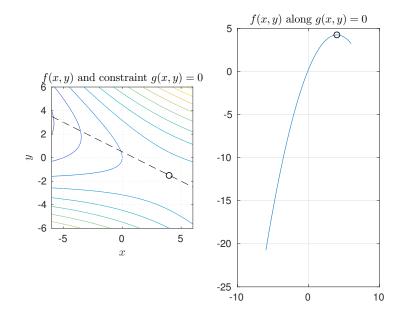
$$2 + y + \lambda = 0$$

$$2y + x + 2\lambda = 0$$

$$x + 2y - 1 = 0$$

where the last is just the constraint.

Step 2: solving these equations tells us that the solution is at:



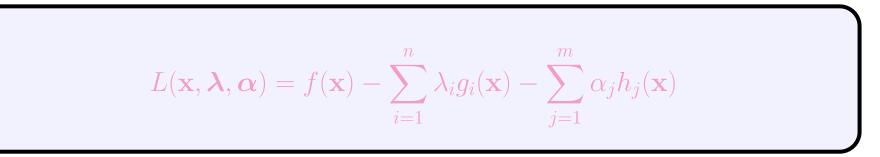
$$(x,y) = (4, -\frac{3}{2})$$

With multiple constraints we follow the same approach, with a *Lagrange multiplier for each constraint*.

How about the full problem? Find

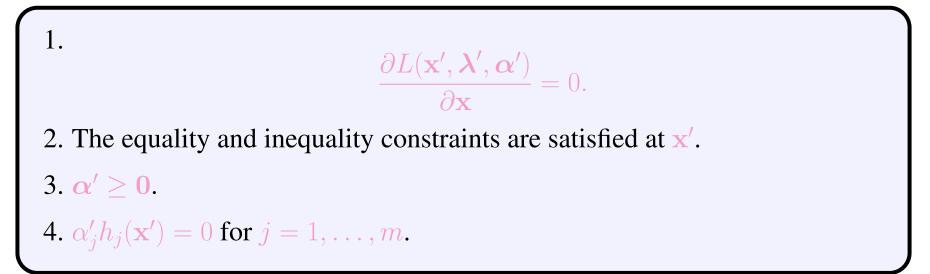
$$\mathbf{x}_{opt} = \operatorname*{argmin}_{\mathbf{x}} f(\mathbf{x})$$
 such that $g_i(\mathbf{x}) = 0$ for $i = 1, 2, \dots, n$
 $h_j(\mathbf{x}) \ge 0$ for $j = 1, 2, \dots, m$

The Lagrangian is now



and the relevant necessary conditions are more numerous.

The necessary conditions now require that when x' is a solution $\exists \lambda', \alpha'$ such that



These are called the *Karush-Kuhn-Tucker (KKT) conditions*.

The *KKT conditions* tell us some important things about the solution. We will only need to address this problem when the constraints are *all inequalities*.

What we've seem so far is called the *primal problem*.

There is also a *dual* version of the problem. Simplifying a little by dropping the equality constraints.

1. The *dual objective function* is

$$\tilde{L}(\boldsymbol{\alpha}) = \inf_{\mathbf{x}} L(\mathbf{x}, \boldsymbol{\alpha}).$$

2. The *dual optimization problem* is

 $\max_{\alpha} \tilde{L}(\alpha) \text{ such that } \alpha \geq \mathbf{0}.$

Sometimes it is *easier to work by solving the dual problem* and this allows us to obtain actual learning algorithms.

We won't be looking in detail at methods for solving such problems, only the *minimum needed to see how SVMs work*.

For the full story see *Numerical Optimization*, Jorge Nocedal and Stephen J. Wright, Second Edition, Springer 2006.

It turns out that with SVMs we get particular benefits when using the *kernel trick*. So we work, as before, in the *extended space*, but now with:

$$f_{\mathbf{w},w_0}(\mathbf{x}) = w_0 + \mathbf{w}^T \mathbf{\Phi}(\mathbf{x})$$
$$h_{\mathbf{w},w_0}(\mathbf{x}) = \operatorname{sgn} \left(f_{\mathbf{w},w_0}(\mathbf{x}) \right)$$

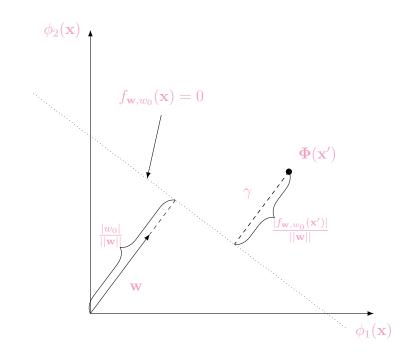
where

$$\operatorname{sgn}(z) = \begin{cases} +1 & \text{if } z > 0 \\ -1 & \text{otherwise.} \end{cases}$$

Note the following:

- 1. Things are easier for SVMs if we use labels $\{+1, -1\}$ for the two classes. (Previously we used $\{0, 1\}$.)
- 2. It also turns out to be easier if we keep w_0 separate rather than rolling it into w.
- 3. We now classify using a "hard" threshold sgn, rather than the "soft" threshold σ .

Consider the geometry again. *Step 1:*



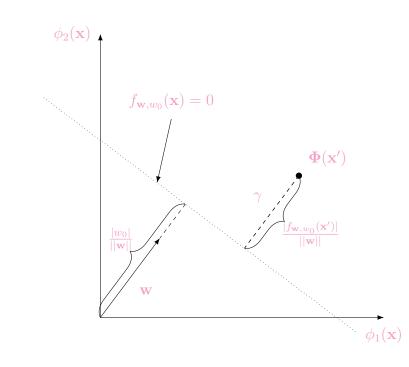
1. We're classifying using the sign of the function

$$f_{\mathbf{w},w_0}(\mathbf{x}) = w_0 + \mathbf{w}^T \mathbf{\Phi}(\mathbf{x}).$$

2. The distance from any point $\Phi(\mathbf{x}')$ in the extended space to the line is



Step 2:



- But we also want the examples to fall on the correct *side* of the line according to their *label*.
- Noting that for any labelled example (\mathbf{x}_i, y_i) the quantity $y_i f_{\mathbf{w}, w_0}(\mathbf{x}_i)$ will be positive if the resulting classification is correct...
- ... the aim is to solve:

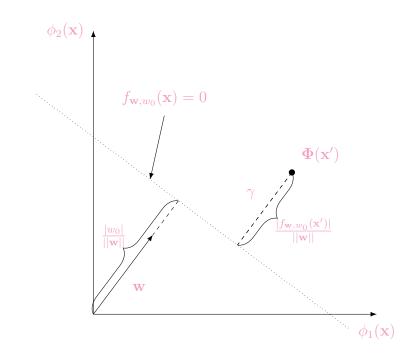
$$(\mathbf{w}, w_o) = \operatorname*{argmax}_{\mathbf{w}, w_0} \left[\min_i \frac{y_i f_{\mathbf{w}, w_0}(\mathbf{x}_i)}{||\mathbf{w}||} \right].$$

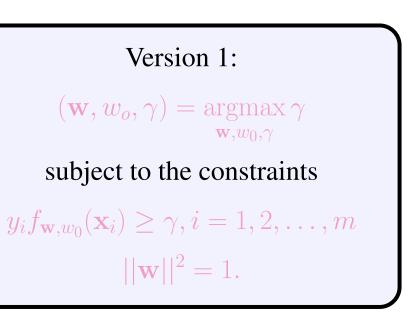
YUK!!!

(With bells on...)

Solution, version 1: convert to a constrained optimization. For any $c \in \mathbb{R}$ $f_{\mathbf{w},w_0}(\mathbf{x}) = 0 \iff w_0 + \mathbf{w}^T \mathbf{\Phi}(\mathbf{x}) = 0$ $\iff cw_0 + c\mathbf{w}^T \mathbf{\Phi}(\mathbf{x}) = 0.$

That means you can fix $||\mathbf{w}||$ to be *anything you like*! (Actually, fix $||\mathbf{w}||^2$ to avoid a square root.)

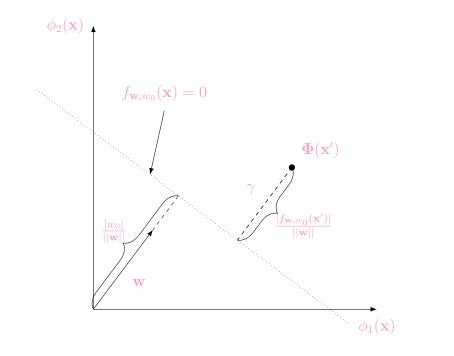


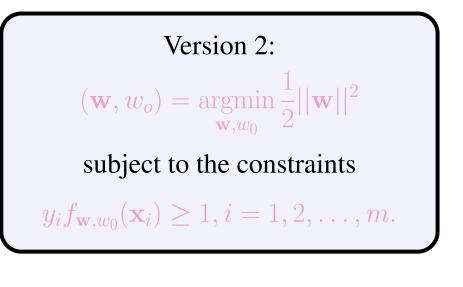


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Solution, version 2: still, convert to a constrained optimization, but instead of fixing $||\mathbf{w}||$:

Fix $\min\{y_i f_{\mathbf{w},w_0}(\mathbf{x}_i)\}$ to be *anything you like*!





(This works because maximizing γ now corresponds to *minimizing* $||\mathbf{w}||$.)

We'll use the second formulation. (You can work through the first as an *exercise*.) The *constrained optimization problem* is:

> Minimize $\frac{1}{2}||\mathbf{w}||^2$ such that $y_i f_{\mathbf{w},w_0}(\mathbf{x}_i) \ge 1$ for $i = 1, \dots, m$.

Referring back, this means the *Lagrangian* is

$$L(\mathbf{w}, w_0, \boldsymbol{\alpha}) = \frac{1}{2} ||\mathbf{w}||^2 - \sum_{i=1}^m \alpha_i (y_i f_{\mathbf{w}, w_0}(\mathbf{x}_i) - 1)$$

and a *necessary condition* for a solution is that

$$\frac{\partial L(\mathbf{w}, w_0, \boldsymbol{\alpha})}{\partial \mathbf{w}} = 0 \qquad \qquad \frac{\partial L(\mathbf{w}, w_0, \boldsymbol{\alpha})}{\partial w_0} = 0.$$

Working these out is easy:

$$\begin{aligned} \frac{\partial L(\mathbf{w}, w_0, \boldsymbol{\alpha})}{\partial \mathbf{w}} &= \frac{\partial}{\partial \mathbf{w}} \left(\frac{1}{2} ||\mathbf{w}||^2 - \sum_{i=1}^m \alpha_i \left(y_i f_{\mathbf{w}, w_0}(\mathbf{x}_i) - 1 \right) \right) \\ &= \mathbf{w} - \sum_{i=1}^m \alpha_i y_i \frac{\partial}{\partial \mathbf{w}} \left(\mathbf{w}^T \mathbf{\Phi}(\mathbf{x}_i) + w_0 \right) \\ &= \mathbf{w} - \sum_{i=1}^m \alpha_i y_i \mathbf{\Phi}(\mathbf{x}_i) \end{aligned}$$

and

$$egin{aligned} rac{\partial L(\mathbf{w},w_0,m{lpha})}{\partial w_0} &= -rac{\partial}{\partial w_0} \left(\sum_{i=1}^m lpha_i y_i f_{\mathbf{w},w_0}(\mathbf{x}_i)
ight) \ &= -rac{\partial}{\partial w_0} \left(\sum_{i=1}^m lpha_i y_i \left(\mathbf{w}^T \mathbf{\Phi}(\mathbf{x}_i) + w_0
ight)
ight) \ &= -\sum_{i=1}^m lpha_i y_i. \end{aligned}$$

Equating those to 0 and adding the *KKT conditions* tells us several things:

1. The weight vector can be expressed as

$$\mathbf{w} = \sum_{i=1}^{m} \alpha_i y_i \mathbf{\Phi}(\mathbf{x}_i)$$

with $\alpha \geq 0$. This is important: we'll return to it in a moment.

2. There is a constraint that

$$\sum_{i=1}^{m} \alpha_i y_i = 0.$$

This will be needed for working out the *dual Lagrangian*.

3. For each example

 $\alpha_i[y_i f_{\mathbf{w},w_0}(\mathbf{x}_i) - 1] = 0.$

The fact that for each example

 $\alpha_i[y_i f_{\mathbf{w},w_0}(\mathbf{x}_i) - 1] = 0$

means that:

Either $y_i f_{\mathbf{w},w_0}(\mathbf{x}_i) = 1$ or $\alpha_i = 0$.

This means that examples fall into two groups.

1. Those for which $y_i f_{\mathbf{w},w_0}(\mathbf{x}_i) = 1$.

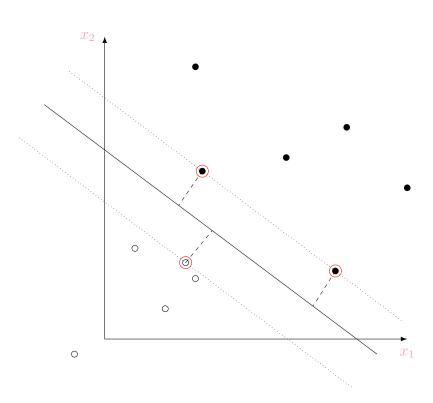
As the contraint used to maximize the margin was $y_i f_{\mathbf{w},w_0}(\mathbf{x}_i) \ge 1$ these are *the examples that are closest to the boundary.*

They are called *support vectors* and they can have *non-zero weights*.

2. Those for which $y_i f_{\mathbf{w},w_0}(\mathbf{x}_i) \neq 1$.

These are non-support vectors and in this case it must be that $\alpha_i = 0$.

Support vectors:



- 1. *Circled examples:* support vectors with $\alpha_i > 0$.
- 2. Other examples: have $\alpha_i = 0$.

Remember that

$$\mathbf{w} = \sum_{i=1}^{m} \alpha_i y_i \mathbf{\Phi}(\mathbf{x}_i).$$

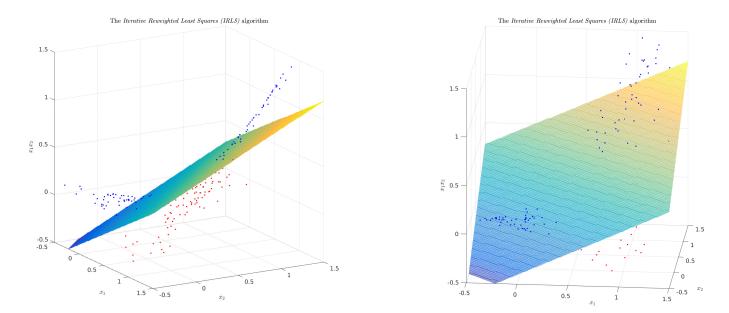
so the weight vector w only depends on the support vectors.

ALSO: the dual parameters α can be used as an *alternative* set of weights. The overall classifier is

$$egin{aligned} h_{\mathbf{w},w_0}(\mathbf{x}) &= \mathrm{sgn}\left(w_0 + \mathbf{w}^T \mathbf{\Phi}(\mathbf{x})
ight) \ &= \mathrm{sgn}\left(w_0 + \sum_{i=1}^m lpha_i y_i \mathbf{\Phi}^T(\mathbf{x}_i) \mathbf{\Phi}(\mathbf{x})
ight) \ &= \mathrm{sgn}\left(w_0 + \sum_{i=1}^m lpha_i y_i K(\mathbf{x}_i,\mathbf{x})
ight) \end{aligned}$$

where $K(\mathbf{x}_i, \mathbf{x}) = \mathbf{\Phi}^T(\mathbf{x}_i) \mathbf{\Phi}(\mathbf{x})$ is called the *kernel*.

Remember where this process started:



The kernel is computing

$$K(\mathbf{x}, \mathbf{x}') = \mathbf{\Phi}^T(\mathbf{x})\mathbf{\Phi}(\mathbf{x}')$$
$$= \sum_{i=1}^k \phi_i(\mathbf{x})\phi_i(\mathbf{x}')$$

This is generally called an *inner product*.

If it's a *hard problem* then you'll probably want *lots of basis functions* so *k is BIG*:

$$egin{aligned} h_{\mathbf{w},w_0}(\mathbf{x}) &= \mathrm{sgn}\left(w_0 + \mathbf{w}^T \mathbf{\Phi}(\mathbf{x})
ight) \ &= \mathrm{sgn}\left(w_0 + \sum_{i=1}^k w_i \phi_i(\mathbf{x})
ight) \ &= \mathrm{sgn}\left(w_0 + \sum_{i=1}^m lpha_i y_i \mathbf{\Phi}^T(\mathbf{x}_i) \mathbf{\Phi}(\mathbf{x})
ight) \ &= \mathrm{sgn}\left(w_0 + \sum_{i=1}^m lpha_i y_i K(\mathbf{x}_i,\mathbf{x})
ight) \end{aligned}$$

What if $K(\mathbf{x}, \mathbf{x}')$ is easy to compute even if k is *HUGE*? (In particular k >> m.)

- 1. We get a definite computational advantage by using the dual version with weights α .
- 2. *Mercer's theorem* tells us exactly when a function K has a corresponding set of *basis functions* $\{\phi_i\}$.

Designing good kernels K is a subject in itself.

Luckily for the majority of the time you will tend to see one of the following:

1. Polynomial:

$$K_{c,d}(\mathbf{x}, \mathbf{x}') = (c + \mathbf{x}^T \mathbf{x}')^d$$

where c and d are parameters.

2. Radial basis function (RBF):

$$K_{\sigma^2}(\mathbf{x}, \mathbf{x}') = \exp\left(-\frac{1}{2\sigma^2}||\mathbf{x} - \mathbf{x}'||^2\right)$$

where σ^2 is a parameter.

The last is particularly prominent. Interestingly, the corresponding set of basis functions is *infinite*. (So we get an improvement in computational complexity from infinite to *linear in the number of examples*!)

Collecting together some of the results up to now:

1. The Lagrangian is

$$L(\mathbf{w}, w_0, \boldsymbol{\alpha}) = \frac{1}{2} ||\mathbf{w}||^2 - \sum_i \alpha_i (y_i f_{\mathbf{w}, w_0}(\mathbf{x}_i) - 1).$$

2. The weight vector is

$$\mathbf{w} = \sum_{i} \alpha_{i} y_{i} \mathbf{\Phi}(\mathbf{x}_{i}).$$

3. The KKT conditions require

$$\sum_{i} \alpha_i y_i = 0.$$

It's easy to show (this is an *exercise*) that the *dual optimization problem* is to maximize

$$\tilde{L}(\boldsymbol{\alpha}) = \sum_{i} \alpha_{i} - \frac{1}{2} \sum_{i} \sum_{j} \alpha_{i} \alpha_{j} y_{i} y_{j} K(\mathbf{x}_{i}, \mathbf{x}_{j})$$

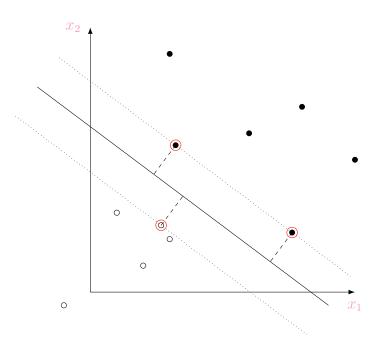
such that $\alpha \geq 0$.

There is one thing still missing:

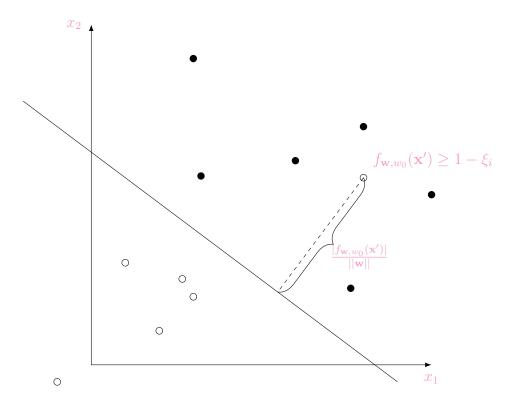
Problem: so far we've only covered the *linearly separable* case.

Even though that means linearly separable *in the extended space* it's still not enough.

By dealing with this we get the *Support Vector Machine (SVM)*.



Fortunately a small modification allows us to let *some* examples be misclassified.



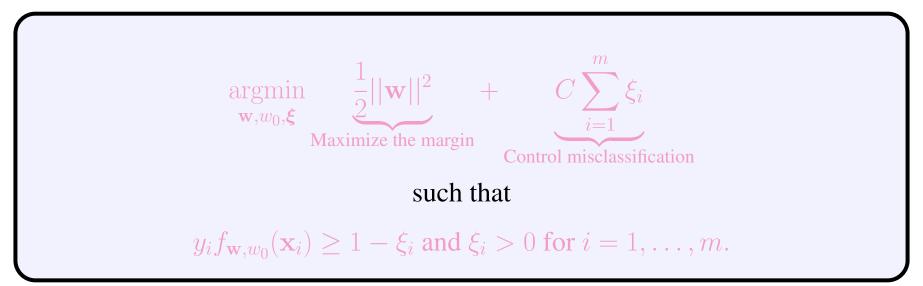
We introduce the *slack variables* ξ_i , one for *each example*.

Although $f_{\mathbf{w},w_0}(\mathbf{x}') < 0$ we have $f_{\mathbf{w},w_0}(\mathbf{x}') \ge 1 - \xi_i$ and we try to force ξ_i to be small.

The constrained optimization problem was:

 $\operatorname{argmin}_{\mathbf{w},w_0} \frac{1}{2} ||\mathbf{w}||^2$ such that $y_i f_{\mathbf{w},w_0}(\mathbf{x}_i) \ge 1$ for $i = 1, \ldots, m$.

The *constrained optimization problem* is now modified to:



There is a *further new parameter* C that controls the trade-off between *maximizing the margin* and *controlling misclassification*.

Support Vector Machines

Once again, the theory of *constrained optimization* can be employed:

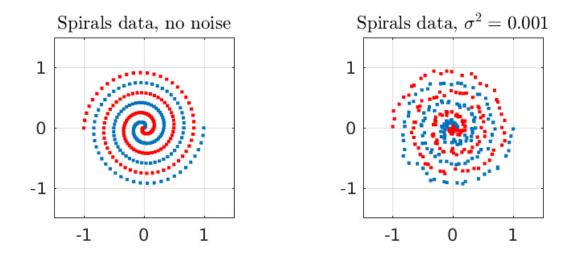
- 1. We get the *same insights* into the solution of the problem, and the *same conclusions*.
- 2. The development is exactly analogous to what we've just seen.

However as is often the case it is not straightforward to move all the way to having a functioning training algorithm.

For this some attention to good *numerical computing* is required. See:

Fast training of support vector machine using sequential minimal optimization, J. C. Platt, *Advances in Kernel Methods*, MIT Press 1999.

Support Vector Machines

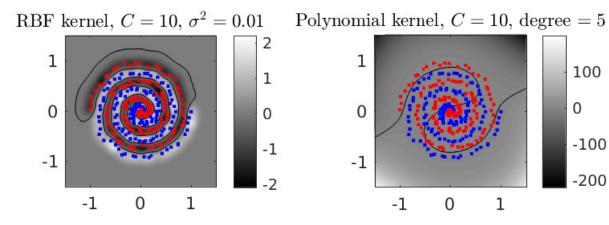


100

-100

-200

0



Supervised learning in practice

We now look at several issues that need to be considered when *applying machine learning algorithms in practice*:

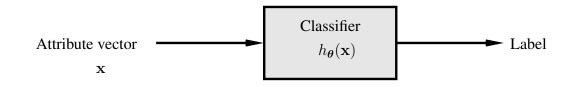
- We often have more examples from some classes than from others.
- The *obvious* measure of performance is not always the *best*.
- Much as we'd love to have an optimal method for *finding hyperparameters*, we don't have one, and it's *unlikely that we ever will*.
- We need to exercise care if we want to claim that one approach is superior to another.

This part of the course has an unusually large number of Commandments.

That's because so many people get so much of it wrong!.

Supervised learning

As usual, we want to design a *classifier*.



It should take an attribute vector

$$\mathbf{x}^T = \begin{bmatrix} x_1 \ x_2 \ \cdots \ x_n \end{bmatrix}$$

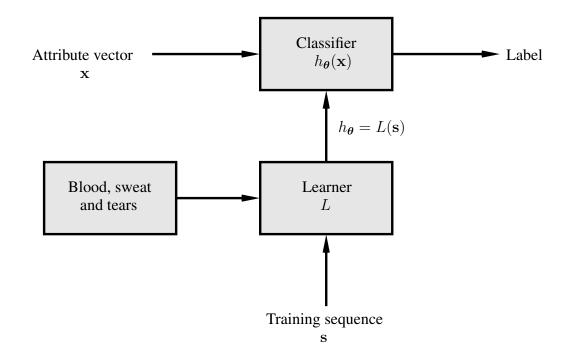
and label it.

We now denote a classifier by $h_{\theta}(\mathbf{x})$ where $\theta^T = (\mathbf{w} \mathbf{p})$ denotes any weights \mathbf{w} and (hyper)parameters \mathbf{p} .

To keep the discussion and notation simple we assume a *classification problem* with *two classes* labelled +1 (*positive examples*) and -1 (*negative examples*).

Supervised learning

Previously, the learning algorithm was a box labelled L.



Unfortunately that turns out not to be enough, so *a new box has been added*.

We've already come across the Commandment:

Thou shalt *try a simple method*. Preferably *many* simple methods.

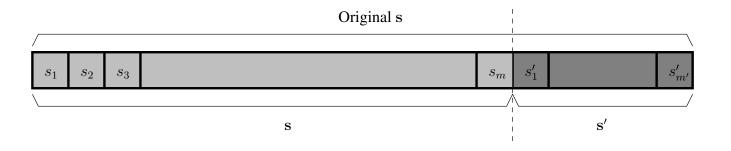
Now we will add:

Thou shalt use an appropriate measure of performance.

Measuring performance

How do you assess the performance of your classifier?

- 1. That is, *after training*, how do you know how well you've done?
- 2. In general, the only way to do this is to divide your examples into a smaller *training set* s of m examples and a *test set* s' of m' examples.



The GOLDEN RULE: data used to assess performance must NEVER have been seen during training.

This might seem obvious, but it was a major flaw in a lot of early work.

How do we choose m and m'? Trial and error!

Assume the training is complete, and we have a classifier h_{θ} obtained using only s. How do we use s' to assess our method's performance?

The obvious way is to see how many examples in s' the classifier classifies correctly:

$$\hat{\operatorname{er}}_{\mathbf{s}'}(h_{\boldsymbol{\theta}}) = \frac{1}{m'} \sum_{i=1}^{m'} \mathbb{I}\left[h_{\boldsymbol{\theta}}(\mathbf{x}'_i) \neq y'_i\right]$$

where

$$\mathbf{s}' = \begin{bmatrix} (\mathbf{x}'_1, y'_1) & (\mathbf{x}'_2, y'_2) & \cdots & (\mathbf{x}'_{m'}, y'_{m'}) \end{bmatrix}^T$$

and

$$\mathbb{I}[z] = \begin{cases} 1 & \text{if } z = \text{true} \\ 0 & \text{if } z = \text{false} \end{cases}.$$

This is just an estimate of the *probability of error* and is often called the *accuracy*.

Unbalanced data

Unfortunately it is often the case that we have *unbalanced data* and this can make such a measure misleading. For example:

If the data is naturally such that *almost all examples are negative* (medical diagnosis for instance) then simply *classifying everything as negative* gives a high performance using this measure.

We need more subtle measures.

For a classifier h and any set s of size m containing m^+ positive examples and m^- negative examples...

Unbalanced data

Define

1. The true positives

$$P^+ = \{(\mathbf{x}, +1) \in \mathbf{s} | h(\mathbf{x}) = +1\}, \text{ and } p^+ = |P^+|$$

2. The *false positives*

$$P^- = \{(\mathbf{x}, -1) \in \mathbf{s} | h(\mathbf{x}) = +1\}, \text{ and } p^- = |P^-|$$

3. The *true negatives*

$$N^+ = \{(\mathbf{x}, -1) \in \mathbf{s} | h(\mathbf{x}) = -1\}, \text{ and } n^+ = |N^+|$$

4. The *false negatives*

$$N^- = \{(\mathbf{x}, +1) \in \mathbf{s} | h(\mathbf{x}) = -1\}, \text{ and } n^- = |N^-|$$

Thus $\hat{er}_{s}(h) = (p^{+} + n^{+})/m$.

This allows us to define more discriminating measures of performance.

Some standard performance measures:

1. Precision $\frac{p^+}{p^++p^-}$. 2. Recall $\frac{p^+}{p^+ + n^-}$. 3. Sensitivity $\frac{p^+}{p^++n^-}$. 4. Specificity $\frac{n^+}{n^++n^-}$. 5. False positive rate $\frac{p^-}{p^-+n^+}$. 6. Positive predictive value $\frac{p^+}{p^++p^-}$. 7. Negative predictive value $\frac{n^+}{n^++n^-}$. 8. False discovery rate $\frac{p^-}{p^-+p^+}$.

In addition, plotting sensitivity (true positive rate) against the false positive rate while a parameter is varied gives the *receiver operating characteristic (ROC)* curve.

Performance measures

The following specifically take account of unbalanced data:

1. Matthews Correlation Coefficient (MCC)

$$\mathbf{MCC} = \frac{p^+ n^+ - p^- n^-}{\sqrt{(p^+ + p^-)(n^+ + n^-)(p^+ + n^-)(n^+ + p^-)}}$$

2. F1 score

$$F1 = \frac{2 \times \text{precision} \times \text{recall}}{\text{precision} + \text{recall}}$$

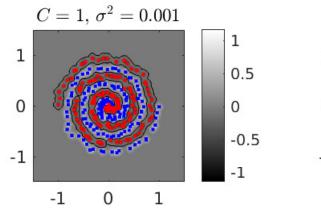
When data is unbalanced these are preferred over the accuracy.

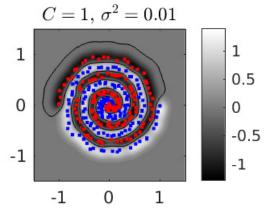
Machine Learning Commandments

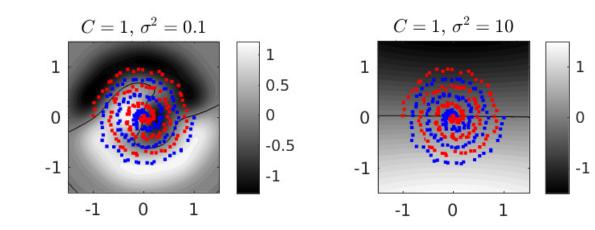
Thou shalt not use *default parameters*.

Thou shalt not use parameters chosen by an *unprincipled formula*. Thou shalt not avoid this issue by clicking on 'Learn' and *hoping it works*. Thou shalt either *choose them carefully* or *integrate them out*.

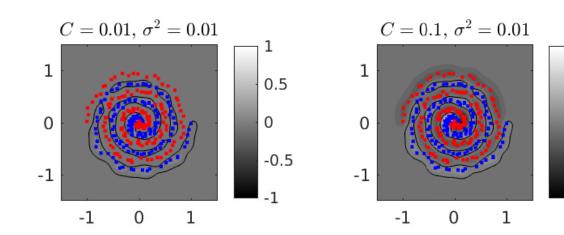
Bad hyperparameters give bad performance

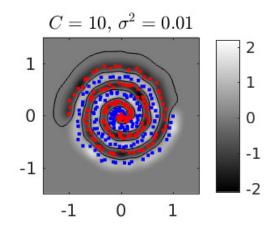


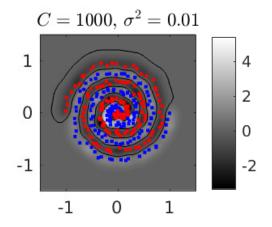




Bad hyperparameters give bad performance







1

0.5

0

-0.5

-1

Validation and crossvalidation

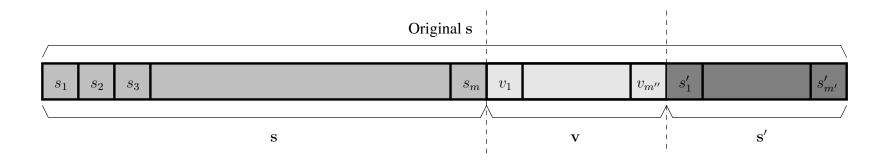
The next question: how do we choose hyperparameters?

Answer: try different values and see which values give the best (estimated) performance.

There is however a problem:

If I use my test set s' to find good hyperparameters, *then I can't use it to get a final measure of performance*. (See the Golden Rule above.)

Solution 1: make a further division of the complete set of examples to obtain a third, *validation* set:



Validation and crossvalidation

Now, to choose the value of a hyperparameter *p*:

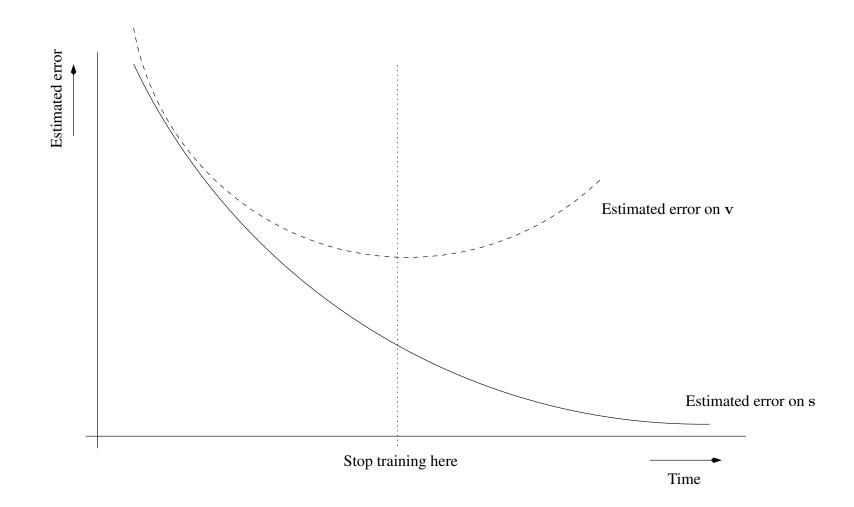
For some range of values p_1, p_2, \ldots, p_n

- 1. Run the training algorithm using training data s and with the hyperparameter set to p_i .
- 2. Assess the resulting h_{θ} by computing a suitable measure (for example accuracy, MCC or F1) using v.

Finally, select the h_{θ} with maximum estimated performance and assess its *actual* performance using s'.

Validation and crossvalidation

This was originally used in a similar way when deciding the best point at which to *stop training* a neural network.

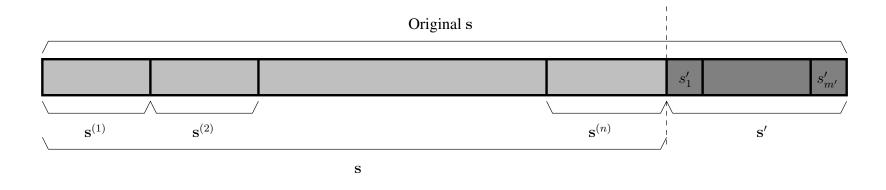


The figure shows the typical scenario.

The method of *crossvalidation* takes this a step further.

We our complete set into training set s and testing set s' as before.

But now instead of further subdividing s just once we divide it into n folds s⁽ⁱ⁾ each having m/n examples.



Typically n = 10 although other values are also used, for example if n = m we have *leave-one-out* cross-validation.

Let s_{-i} denote the set obtained from s by *removing* $s^{(i)}$.

Let $\hat{er}_{s(i)}(h)$ denote any suitable error measure, such as accuracy, MCC or F1, computed for *h* using fold *i*.

Let $L_{s_{-i},p}$ be the classifier obtained by running learning algorithm L on examples s_{-i} using hyperparameters p.

Then,

$$\frac{1}{n}\sum_{i=1}^{n} \hat{\operatorname{er}}_{\mathbf{s}^{(i)}}(L_{\mathbf{s}_{-i},\mathbf{p}})$$

is the *n*-fold crossvalidation error estimate.

So for example, let $s_j^{(i)}$ denote the *j*th example in the *i*th fold. Then using accuracy as the error estimate we have

$$\frac{1}{m} \sum_{i=1}^{n} \sum_{j=1}^{m/n} \mathbb{I}\left[L_{\mathbf{s}_{-i},\mathbf{p}}(\mathbf{x}_{j}^{(i)}) \neq y_{j}^{(i)}\right]$$

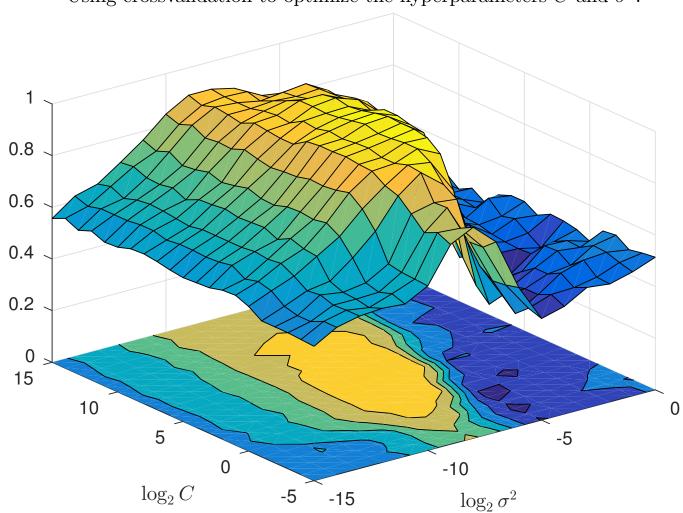
Two further points:

- 1. What if the data are unbalanced? *Stratified crossvalidation* chooses folds such that the proportion of positive examples in each fold matches that in s.
- 2. Hyperparameter choice can be done just as above, using a basic search.

What happens however if we have multiple hyperparameters?

- 1. We can search over all combinations of values for specified ranges of each parameter.
- 2. This is the standard method in choosing parameters for support vector machines (SVMs).
- 3. With SVMs it is generally limited to the case of only two hyperparameters.
- 4. Larger numbers quickly become infeasible.

This is what we get for an *SVM* applied to the *two spirals*:



Using crossvalidation to optimize the hyperparameters C and σ^2 .

Machine Learning Commandments

Thou shalt *provide evidence* before claiming that *thy method is the best*.

The shalt take extra notice of this Commandment if *thou considers thyself a True And Pure Bayesian*.

Imagine I have compared the *Bloggs Classificator 2000* and the *CleverCorp Discriminotron* and found that:

1. Bloggs Classificator 2000 has estimated accuracy 0.981 on the test set.

2. CleverCorp Discriminotron has estimated accuracy 0.982 on the test set.

Can I claim that the CleverCorp Discriminotron is the better classifier? Answer:

NO!!!!!!!

Note for next year: include photo of grumpy-looking cat.

Assessing a single classifier

From Mathematical Methods for Computer Science:

The *Central Limit Theorem*: If we have independent identically distributed (iid) random variables X_1, X_2, \ldots, X_n with mean

$$\mathbb{E}\left[X\right] = \mu$$

and standard deviation

$$\mathbb{E}\left[(X-\mu)^2\right] = \sigma^2$$

then as $n \to \infty$

$$\frac{\hat{X}_n - \mu}{\sigma/\sqrt{n}} \to N(0, 1)$$

where

$$\hat{X}_n = \frac{1}{n} \sum_{i=1}^n X_i.$$

We have tables of values z_p such that if $x \sim N(0, 1)$ then

 $\Pr\left(-z_p \le x \le z_p\right) > p.$

Rearranging this using the equation from the previous slide we have that with probability p

$$\mu \in \left[\hat{X}_n \pm z_p \sqrt{\frac{\sigma^2}{n}}\right].$$

We don't know σ^2 but it can be estimated using

$$\sigma^2 \simeq \frac{1}{n-1} \sum_{i=1}^n \left(X_i - \hat{X}_n \right)^2.$$

Alternatively, when X takes only values 0 or 1

$$\sigma^2 = \mathbb{E}\left[(X - \mu)^2 \right] = \mathbb{E}\left[X^2 \right] - \mu^2 = \mu(1 - \mu) \simeq \hat{X}_n(1 - \hat{X}_n).$$

The *actual probability of error* for a classifier h is $er(h) = \mathbb{E} \left[\mathbb{I} \left[h(\mathbf{x}) \neq y \right] \right]$

and we are $\textit{estimating}\; \texttt{er}(h)$ using the accuracy

$$\hat{\operatorname{er}}_{\mathbf{s}}(h) = \frac{1}{m} \sum_{i=1}^{m} \mathbb{I}[h(\mathbf{x}_i) \neq y_i]$$

for a test set s.

We can find a confidence interval for this estimate using precisely the derivation above, simply by noting that the X_i are the random variables

 $X_i = \mathbb{I}\left[h(\mathbf{x}_i) \neq y_i\right].$

Typically we are interested in a 95% confidence interval, for which $z_p = 1.96$.

Thus, when m > 30 (so that the central limit theorem applies) we know that, with probability 0.95

$$\mathbf{er}(h) = \hat{\mathbf{er}}_{\mathbf{s}}(h) \pm 1.96\sqrt{\frac{\hat{\mathbf{er}}_{\mathbf{s}}(h)(1 - \hat{\mathbf{er}}_{\mathbf{s}}(h)))}{m}}$$

Example: I have 100 test examples and my classifier makes 18 errors. With probability 0.95 I know that

$$er(h) = 0.18 \pm 1.96 \sqrt{\frac{0.18(1 - 0.18)}{100}}$$
$$= 0.18 \pm 0.075.$$

This should perhaps *raise an alarm* regarding our suggested comparison of classifiers above.

There is an important distinction to be made here:

- 1. The *mean of* X is μ and the *variance of* X is σ^2 .
- 2. We can also ask about the mean and variance of \hat{X}_n .
- 3. The *mean of* \hat{X}_n is

$$\mathbb{E}\left[\hat{X}_{n}\right] = \mathbb{E}\left[\frac{1}{n}\sum_{i=1}^{n}X_{i}\right]$$
$$= \frac{1}{n}\sum_{i=1}^{n}\mathbb{E}\left[X_{i}\right]$$
$$= \mu.$$

4. It is left as an *exercise* to show that the *variance of* \hat{X}_n is

$$\sigma_{\hat{X}_n}^2 = \frac{\sigma^2}{n}$$

We are using the values z_p such that if $x \sim N(0, 1)$ then

 $\Pr(-z_p \le x \le z_p) > p.$

There is an *alternative* way to think about this.

- 1. Say we have a random variable Y with variance σ_Y^2 and mean μ_Y .
- 2. The random variable $Y \mu_Y$ has variance σ_Y^2 and mean 0.
- 3. It is a straightforward exercise to show that dividing a random variable having variance σ^2 by σ gives us a new random variable with variance 1.
- 4. Thus the random variable $\frac{Y-\mu_Y}{\sigma_Y}$ has mean 0 and variance 1.

So: with probability *p*

$$Y = \mu_Y \pm z_p \sigma_Y$$
$$\mu_Y = Y \pm z_p \sigma_Y.$$

Compare this with what we saw earlier. You need to be careful to keep track of whether you are considering the mean and variance of a single RV or a sum of RVs.

Now say I have classifiers h_1 (*Bloggs Classificator 2000*) and h_2 (*CleverCorp Discriminotron*) and I want to know something about the quantity

 $d = \operatorname{er}(h_1) - \operatorname{er}(h_2).$

I estimate d using

 $\hat{d} = \hat{\operatorname{er}}_{\mathbf{s}_1}(h_1) - \hat{\operatorname{er}}_{\mathbf{s}_2}(h_2)$

where s_1 and s_2 are *two* independent test sets.

Notice:

- 1. The estimate of *d* is a sum of random variables, and *we can apply the central limit theorem*.
- 2. The estimate is *unbiased*

 $\mathbb{E}\left[\hat{\operatorname{er}}_{\mathbf{s}_1}(h_1) - \hat{\operatorname{er}}_{\mathbf{s}_2}(h_2)\right] = d.$

Also notice:

- 1. The two parts of the estimate $\hat{er}_{s_1}(h_1)$ and $\hat{er}_{s_2}(h_2)$ are each sums of random variables and *we can apply the central limit theorem to each*.
- 2. The variance of the estimate is the sum of the variances of $\hat{er}_{s_1}(h_1)$ and $\hat{er}_{s_2}(h_2)$.
- 3. Adding Gaussians gives another Gaussian.
- 4. We can calculate a confidence interval for our estimate.

With probability 0.95

$$d = \hat{d} \pm 1.96 \sqrt{\frac{\hat{\mathrm{er}}_{\mathbf{s}_1}(h_1)(1 - \hat{\mathrm{er}}_{\mathbf{s}_1}(h_1))}{m_1} + \frac{\hat{\mathrm{er}}_{\mathbf{s}_2}(h_2)(1 - \hat{\mathrm{er}}_{\mathbf{s}_2}(h_2))}{m_2}}$$

In fact, if we are using a split into training set s and test set s' we can generally obtain h_1 and h_2 using s and use the estimate

$$\hat{d} = \hat{\operatorname{er}}_{\mathbf{s}'}(h_1) - \hat{\operatorname{er}}_{\mathbf{s}'}(h_2).$$

Comparing classifiers—hypothesis testing

This still doesn't tell us directly about *whether one classifier is better than another*—whether h_1 is better than h_2 .

What we actually want to know is whether

 $d = \operatorname{er}(h_1) - \operatorname{er}(h_2) > 0.$

Say we've measured $\hat{D} = \hat{d}$. Then:

- Imagine the *actual value* of *d* is 0.
- Recall that the *mean* of \hat{D} is d.
- So *larger* measured values \hat{d} are *less likely*, even though some random variation is inevitable.
- If it is highly *unlikely* that when d = 0 a measured value of \hat{d} would be observed, then we can be confident that d > 0.
- Thus we are interested in

$$\Pr(\hat{D} > d + \hat{d}).$$

This is known as a *one-sided bound*.

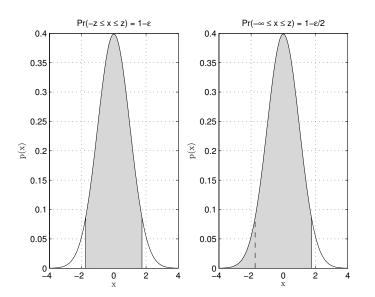
One-sided bounds

Given the *two-sided bound*

 $\Pr(-z_{\epsilon} \le x \le z_{\epsilon}) = 1 - \epsilon$

we actually need to know the *one-sided bound*

 $\Pr(x \le z_{\epsilon}).$



Clearly, if our random variable is *Gaussian* then $\Pr(x \le z_{\epsilon}) = 1 - \epsilon/2$.

We now know how to compare *hypotheses* h_1 and h_2 .

But we still haven't properly addressed the comparison of *algorithms*.

- Remember, a learning algorithm L maps training data s to hypothesis h.
- So we *really* want to know about the quantity

 $d = \mathbb{E}_{\mathbf{s} \in S^m} \left[\operatorname{er}(L_1(\mathbf{s})) - \operatorname{er}(L_2(\mathbf{s})) \right].$

• This is the *expected difference* between the *actual errors* of the *two different* algorithms L_1 and L_2 .

Unfortunately, we have *only one set of data* s available and we *can only estimate* errors er(h)—we don't have access to the *actual quantities*.

We can however use the idea of *crossvalidation*.

Comparing algorithms: paired t-tests

Recall, we subdivide s into n folds s⁽ⁱ⁾ each having m/n examples



and denote by s_{-i} the set obtained from s by *removing* $s^{(i)}$. Then

 $\frac{1}{n}\sum_{i=1}^{n} \hat{\operatorname{er}}_{\mathbf{s}^{(i)}}(L(\mathbf{s}_{-i}))$

is the *n*-fold crossvalidation error estimate. Now we estimate d using

$$\hat{d} = \frac{1}{n} \sum_{i=1}^{n} \left[\hat{\text{er}}_{\mathbf{s}^{(i)}}(L_1(\mathbf{s}_{-i})) - \hat{\text{er}}_{\mathbf{s}^{(i)}}(L_2(\mathbf{s}_{-i})) \right].$$

As usual, there is a *statistical test* allowing us to assess *how likely this estimate is to mislead us*.

We will not consider the derivation in detail. With probability p

$d \in$	$\left[\hat{d} \pm t_{p,n-1}\sigma_{\hat{d}}\right]$	•
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This is analogous to the equations seen above, however:

- The parameter $t_{p,n-1}$ is analogous to z_p .
- The parameter $t_{p,n-1}$ is related to the area under the *Student's t-distribution* whereas z_p is related to the area under the normal distribution.
- The relevant estimate of *standard deviation* is

$$\sigma_{\hat{d}} = \sqrt{\frac{1}{n(n-1)} \sum_{i=1}^{n} \left(d_i - \hat{d} \right)^2}$$

where

$$d_i = \hat{\operatorname{er}}_{\mathbf{s}^{(i)}}(L_1(\mathbf{s}_{-i})) - \hat{\operatorname{er}}_{\mathbf{s}^{(i)}}(L_2(\mathbf{s}_{-i})).$$