# A Course in 



## 9 | Kernel Methods

Linear models are great because they are easy to understand and easy to optimize. They suffer because they can only learn very simple decision boundaries. Neural networks can learn more complex decision boundaries, but lose the nice convexity properties of many linear models.

One way of getting a linear model to behave non-linearly is to transform the input. For instance, by adding feature pairs as additional inputs. Learning a linear model on such a representation is convex, but is computationally prohibitive in all but very low dimensional spaces. You might ask: instead of explicitly expanding the feature space, is it possible to stay with our original data representation and do all the feature blow up implicitly? Surprisingly, the answer is often "yes" and the family of techniques that makes this possible are known as kernel approaches.

### 9.1 From Feature Combinations to Kernels

In Section 4.4, you learned one method for increasing the expressive power of linear models: explode the feature space. For instance, a "quadratic" feature explosion might map a feature vector $x=$ $\left\langle x_{1}, x_{2}, x_{3}, \ldots, x_{D}\right\rangle$ to an expanded version denoted $\phi(x)$ :

$$
\begin{align*}
\phi(\boldsymbol{x})= & \left\langle 1,2 x_{1}, 2 x_{2}, 2 x_{3}, \ldots, 2 x_{D},\right. \\
& x_{1}^{2}, x_{1} x_{2}, x_{1} x_{3}, \ldots, x_{1} x_{D}, \\
& x_{2} x_{1}, x_{2}^{2}, x_{2} x_{3}, \ldots, x_{2} x_{D}, \\
& x_{3} x_{1}, x_{3} x_{2}, x_{3}^{2}, \ldots, x_{2} x_{D}, \\
& \ldots,  \tag{9.1}\\
& \left.x_{D} x_{1}, x_{D} x_{2}, x_{D} x_{3}, \ldots, x_{D}^{2}\right\rangle
\end{align*}
$$

(Note that there are repetitions here, but hopefully most learning algorithms can deal well with redundant features; in particular, the $2 x_{1}$ terms are due to collapsing some repetitions.)

## Learning Objectives:

- Explain how kernels generalize both feature combinations and basis functions.
- Contrast dot products with kernel products.
- Implement kernelized perceptron.
- Derive a kernelized version of regularized least squares regression.
- Implement a kernelized version of the perceptron.
- Derive the dual formulation of the support vector machine.

Dependencies:

You could then train a classifier on this expanded feature space. There are two primary concerns in doing so. The first is computational: if your learning algorithm scales linearly in the number of features, then you've just squared the amount of computation you need to perform; you've also squared the amount of memory you'll need. The second is statistical: if you go by the heuristic that you should have about two examples for every feature, then you will now need quadratically many training examples in order to avoid overfitting.

This chapter is all about dealing with the computational issue. It will turn out in Chapter ?? that you can also deal with the statistical issue: for now, you can just hope that regularization will be sufficient to attenuate overfitting.

The key insight in kernel-based learning is that you can rewrite many linear models in a way that doesn't require you to ever $e x$ plicitly compute $\phi(x)$. To start with, you can think of this purely as a computational "trick" that enables you to use the power of a quadratic feature mapping without actually having to compute and store the mapped vectors. Later, you will see that it's actually quite a bit deeper. Most algorithms we discuss involve a product of the form $w \cdot \phi(x)$, after performing the feature mapping. The goal is to rewrite these algorithms so that they only ever depend on dot products between two examples, say $x$ and $z$; namely, they depend on $\phi(x) \cdot \phi(z)$. To understand why this is helpful, consider the quadratic expansion from above, and the dot-product between two vectors. You get:

$$
\begin{align*}
\phi(x) \cdot \phi(z)= & 1+x_{1} z_{1}+x_{2} z_{2}+\cdots+x_{D} z_{D}+x_{1}^{2} z_{1}^{2}+\cdots+x_{1} x_{D} z_{1} z_{D}+ \\
& \cdots+x_{D} x_{1} z_{D} z_{1}+x_{D} x_{2} z_{D} z_{2}+\cdots+x_{D}^{2} z_{D}^{2}  \tag{9.2}\\
= & 1+2 \sum_{d} x_{d} z_{d}+\sum_{d} \sum_{e} x_{d} x_{e} z_{d} z_{e}  \tag{9.3}\\
= & 1+2 x \cdot z+(\boldsymbol{x} \cdot \boldsymbol{z})^{2}  \tag{9.4}\\
= & (1+\boldsymbol{x} \cdot \boldsymbol{z})^{2} \tag{9.5}
\end{align*}
$$

Thus, you can compute $\phi(x) \cdot \phi(z)$ in exactly the same amount of time as you can compute $x \cdot z$ (plus the time it takes to perform an addition and a multiply, about 0.02 nanoseconds on a circa 2011 processor).

The rest of the practical challenge is to rewrite your algorithms so that they only depend on dot products between examples and not on any explicit weight vectors.

### 9.2 Kernelized Perceptron

Consider the original perceptron algorithm from Chapter 3, repeated in Algorithm 9.2 using linear algebra notation and using feature expansion notation $\phi(x)$. In this algorithm, there are two places

```
Algorithm 28 PerceptronTrain(D, MaxIter)
    \(w \leftarrow \mathbf{0}, b \leftarrow o \quad / /\) initialize weights and bias
    for iter \(=1 \ldots\) MaxIter do
        for all \((x, y) \in \mathbf{D}\) do
            \(a \leftarrow \boldsymbol{w} \cdot \phi(x)+b \quad / /\) compute activation for this example
            if \(y a \leq o\) then
                \(\boldsymbol{w} \leftarrow w+y \phi(x) \quad\) // update weights
                \(b \leftarrow b+y \quad / /\) update bias
            end if
        end for
    end for
    return \(w, b\)
```


## Math Review | Spans and null spaces

reminder: if $\mathcal{U}=\left\{\boldsymbol{u}_{i}\right\}_{i}$ is a set of vectors in $\mathbb{R}^{D}$, then the span of $\mathcal{U}$ is the set of vectors that can be written as linear combinations of $\boldsymbol{u}_{i} \mathrm{~s}$; namely: $\operatorname{span}(\mathcal{U})=\left\{\sum_{i} a_{i} \boldsymbol{u}_{i}: a_{1} \in \mathbb{R}, \ldots, a_{I} \in \mathbb{R}\right\}$. the null space of $\mathcal{U}$ is everything that's left: $\mathbb{R}^{D} \backslash \operatorname{span}(\mathcal{U})$.
TODO pictures
where $\phi(\boldsymbol{x})$ is used explicitly. The first is in computing the activation (line 4) and the second is in updating the weights (line 6). The goal is to remove the explicit dependence of this algorithm on $\phi$ and on the weight vector.

To do so, you can observe that at any point in the algorithm, the weight vector $w$ can be written as a linear combination of expanded training data. In particular, at any point, $w=\sum_{n} \alpha_{n} \phi\left(x_{n}\right)$ for some parameters $\alpha$. Initially, $\boldsymbol{w}=\mathbf{0}$ so choosing $\alpha=0$ yields this. If the first update occurs on the $n$th training example, then the resolution weight vector is simply $y_{n} \phi\left(x_{n}\right)$, which is equivalent to setting $\alpha_{n}=$ $y_{n}$. If the second update occurs on the $m$ th training example, then all you need to do is update $\alpha_{m} \leftarrow \alpha_{m}+y_{m}$. This is true, even if you make multiple passes over the data. This observation leads to the following representer theorem, which states that the weight vector of the perceptron lies in the span of the training data.

Theorem 11 (Perceptron Representer Theorem). During a run of the perceptron algorithm, the weight vector $\boldsymbol{w}$ is always in the span of the (assumed non-empty) training data, $\phi\left(x_{1}\right), \ldots, \phi\left(x_{N}\right)$.

Proof of Theorem 11. By induction. Base case: the span of any nonempty set contains the zero vector, which is the initial weight vector. Inductive case: suppose that the theorem is true before the $k$ th update, and suppose that the $k$ th update happens on example $n$. By the inductive hypothesis, you can write $w=\sum_{i} \alpha_{i} \phi\left(\boldsymbol{x}_{i}\right)$ before

```
Algorithm 29 KernelizedPerceptronTrain(D, MaxIter)
    \(\alpha \leftarrow 0, b \leftarrow o \quad / /\) initialize coefficients and bias
    for iter \(=1 \ldots\) MaxIter do
        for all \(\left(x_{n}, y_{n}\right) \in \mathbf{D}\) do
            \(a \leftarrow \sum_{m} \alpha_{m} \phi\left(x_{m}\right) \cdot \phi\left(x_{n}\right)+b \quad / /\) compute activation for this example
            if \(y_{n} a \leq o\) then
                \(\alpha_{n} \leftarrow \alpha_{n}+y_{n} \quad\) // update coefficients
                    \(b \leftarrow b+y \quad\) I/ update bias
            end if
        end for
    end for
    return \(\alpha, b\)
```

the update. The new weight vector is $\left[\sum_{i} \alpha_{i} \phi\left(x_{i}\right)\right]+y_{n} \phi\left(x_{n}\right)=$ $\sum_{i}\left(\alpha_{i}+y_{n}[i=n]\right) \phi\left(x_{i}\right)$, which is still in the span of the training data.

Now that you know that you can always write $w=\sum_{n} \alpha_{n} \phi\left(x_{n}\right)$ for some $\alpha_{i} \mathrm{~s}$, you can additionall compute the activations (line 4) as:

$$
\begin{align*}
w \cdot \phi(x)+b & =\left(\sum_{n} \alpha_{n} \phi\left(x_{n}\right)\right) \cdot \phi(x)+b \quad \text { definition of } w \\
& =\sum_{n} \alpha_{n}\left[\phi\left(x_{n}\right) \cdot \phi(x)\right]+b \quad \text { dot products are linear } \tag{9.6}
\end{align*}
$$

This now depends only on dot-products between data points, and never explicitly requires a weight vector. You can now rewrite the entire perceptron algorithm so that it never refers explicitly to the weights and only ever depends on pairwise dot products between examples. This is shown in Algorithm 9.2.

The advantage to this "kernelized" algorithm is that you can perform feature expansions like the quadratic feature expansion from the introduction for "free." For example, for exactly the same cost as the quadratic features, you can use a cubic feature map, computed as $\phi \ddot{\boldsymbol{x}}) \phi(\boldsymbol{z})=(1+\boldsymbol{x} \cdot \boldsymbol{z})^{3}$, which corresponds to three-way interactions between variables. (And, in general, you can do so for any polynomial degree $p$ at the same computational complexity.)

### 9.3 Kernelized K-means

For a complete change of pace, consider the K-means algorithm from Section ??. This algorithm is for clustering where there is no notion of "training labels." Instead, you want to partition the data into coherent clusters. For data in $\mathbb{R}^{D}$, it involves randomly initializing $K$-many
cluster means $\boldsymbol{\mu}^{(1)}, \ldots, \boldsymbol{\mu}^{(\mathrm{K})}$. The algorithm then alternates between the following two steps until convergence, with $x$ replaced by $\phi(x)$ since that is the eventual goal:

1. For each example $n$, set cluster label $z_{n}=\arg \min _{k}\left\|\phi\left(\boldsymbol{x}_{n}\right)-\boldsymbol{\mu}^{(\mathrm{k})}\right\|^{2}$.
2. For each cluster $k$, update $\boldsymbol{\mu}^{(k)}=\frac{1}{N_{k}} \sum_{n: z_{n}=k} \phi\left(x_{n}\right)$, where $N_{k}$ is the number of $n$ with $z_{n}=k$.

The question is whether you can perform these steps without explicitly computing $\phi\left(x_{n}\right)$. The representer theorem is more straightforward here than in the perceptron. The mean of a set of data is, almost by definition, in the span of that data (choose the $a_{i} \mathrm{~s}$ all to be equal to $1 / N)$. Thus, so long as you initialize the means in the span of the data, you are guaranteed always to have the means in the span of the data. Given this, you know that you can write each mean as an expansion of the data; say that $\mu^{(\mathrm{k})}=\sum_{n} \alpha_{n}^{(\mathrm{k})} \phi\left(x_{n}\right)$ for some parameters $\alpha_{n}^{(\mathrm{k})}$ (there are $N \times K$-many such parameters).

Given this expansion, in order to execute step (1), you need to compute norms. This can be done as follows:

$$
\begin{array}{rlr}
z_{n} & =\arg \min _{k}\left\|\phi\left(x_{n}\right)-\mu^{(\mathrm{k})}\right\|^{2} & \text { definition of } z_{n} \\
& =\arg \min _{k}\left\|\phi\left(x_{n}\right)-\sum_{m} \alpha_{m}^{(\mathrm{k})} \phi\left(x_{m}\right)\right\|^{2} & \text { definition of } \mu^{(\mathrm{k})} \\
& =\arg \min _{k}\left\|\phi\left(x_{n}\right)\right\|^{2}+\left\|\sum_{m} \alpha_{m}^{(\mathrm{k})} \phi\left(x_{m}\right)\right\|^{2}+\phi\left(x_{n}\right) \cdot\left[\sum_{m} \alpha_{m}^{(\mathrm{k})} \phi\left(x_{m}\right)\right] & \text { expand quadratic term } \\
& =\arg \min _{k} \sum_{m} \sum_{m^{\prime}} \alpha_{m}^{(\mathrm{k})} \alpha_{m^{\prime}}^{(\mathrm{k})} \phi\left(x_{m}\right) \cdot \phi\left(x_{m^{\prime}}\right)+\sum_{m} \alpha_{m}^{(\mathrm{k})} \phi\left(x_{m}\right) \cdot \phi\left(x_{n}\right)+\text { const } & \text { linearity and constant } \tag{9.11}
\end{array}
$$

This computation can replace the assignments in step (1) of K-means. The mean updates are more direct in step (2):

$$
\boldsymbol{\mu}^{(\mathrm{k})}=\frac{1}{N_{k}} \sum_{n: z_{n}=k} \phi\left(x_{n}\right) \quad \Longleftrightarrow \quad \alpha_{n}^{(\mathrm{k})}= \begin{cases}\frac{1}{N_{k}} & \text { if } z_{n}=k  \tag{9.12}\\ 0 & \text { otherwise }\end{cases}
$$

### 9.4 What Makes a Kernel

A kernel is just a form of generalized dot product. You can also think of it as simply shorthand for $\phi(\boldsymbol{x}) \cdot \phi(\boldsymbol{z})$, which is commonly written $K^{\phi}(x, z)$. Or, when $\phi$ is clear from context, simply $K(x, z)$.

This is often refered to as the kernel product between $x$ and $z$ (under the mapping $\phi$ ).

In this view, what you've seen in the preceding two sections is that you can rewrite both the perceptron algorithm and the K-means algorithm so that they only ever depend on kernel products between data points, and never on the actual datapoints themselves. This is a very powerful notion, as it has enabled the development of a large number of non-linear algorithms essentially "for free" (by applying the so-called kernel trick, that you've just seen twice).

This raises an interesting question. If you have rewritten these algorithms so that they only depend on the data through a function $K: \mathcal{X} \times \mathcal{X} \rightarrow \mathbb{R}$, can you stick any function $K$ in these algorithms, or are there some $K$ that are "forbidden?" In one sense, you "could" use any $K$, but the real question is: for what types of functions $K$ do these algorithms retain the properties that we expect them to have (like convergence, optimality, etc.)?

One way to answer this question is to say that $K(\cdot, \cdot)$ is a valid kernel if it corresponds to the inner product between two vectors. That is, $K$ is valid if there exists a function $\phi$ such that $K(x, z)=$ $\phi(\boldsymbol{x}) \cdot \phi(\boldsymbol{z})$. This is a direct definition and it should be clear that if $K$ satisfies this, then the algorithms go through as expected (because this is how we derived them).

You've already seen the general class of polynomial kernels, which have the form:

$$
\begin{equation*}
K_{d}^{(\text {poly })}(x, z)=(1+x \cdot z)^{d} \tag{9.13}
\end{equation*}
$$

where $d$ is a hyperparameter of the kernel. These kernels correspond to polynomial feature expansions.

There is an alternative characterization of a valid kernel function that is more mathematical. It states that $K: \mathcal{X} \times \mathcal{X} \rightarrow \mathbb{R}$ is a kernel if $K$ is positive semi-definite (or, in shorthand, psd). This property is also sometimes called Mercer's condition. In this context, this means the for all functions $f$ that are square integrable (i.e., $\int f(x)^{2} \mathrm{~d} x<\infty$ ), other than the zero function, the following property holds:

$$
\begin{equation*}
\iint f(x) K(x, z) f(z) \mathrm{d} x \mathrm{~d} z>0 \tag{9.14}
\end{equation*}
$$

This likely seems like it came out of nowhere. Unfortunately, the connection is well beyond the scope of this book, but is covered well is external sources. For now, simply take it as a given that this is an equivalent requirement. (For those so inclined, the appendix of this book gives a proof, but it requires a bit of knowledge of function spaces to understand.)

The question is: why is this alternative characterization useful? It is useful because it gives you an alternative way to construct kernel
functions. For instance, using it you can easily prove the following, which would be difficult from the definition of kernels as inner products after feature mappings.

Theorem 12 (Kernel Addition). If $K_{1}$ and $K_{2}$ are kernels, the $K$ defined by $K(x, z)=K_{1}(x, z)+K_{2}(x, z)$ is also a kernel.
Proof of Theorem 12. You need to verify the positive semi-definite property on $K$. You can do this as follows:

$$
\begin{array}{rlrl}
\iint f(x) K(x, z) f(z) \mathrm{d} x \mathrm{~d} z= & \iint f(x)\left[K_{1}(x, z)+K_{2}(x, z)\right] f(z) \mathrm{d} x \mathrm{~d} z & \text { definition of } K \\
= & \iint f(x) K_{1}(x, z) f(z) \mathrm{d} x \mathrm{~d} z & & \\
& +\iint f(x) K_{2}(x, z) f(z) \mathrm{d} x \mathrm{~d} z & & \text { distributive rule } \\
& >0+0 & (9 \cdot 16) & K_{1} \text { and } K_{2} \text { are psd } \tag{9.17}
\end{array}
$$

More generally, any positive linear combination of kernels is still a kernel. Specifically, if $K_{1}, \ldots, K_{M}$ are all kernels, and $\alpha_{1}, \ldots, \alpha_{M} \geq 0$, then $K(x, z)=\sum_{m} \alpha_{m} K_{m}(x, z)$ is also a kernel.

You can also use this property to show that the following Gaussian kernel (also called the RBF kernel) is also psd:

$$
\begin{equation*}
K_{\gamma}^{(\mathrm{RBF})}(x, z)=\exp \left[-\gamma\|x-z\|^{2}\right] \tag{9.18}
\end{equation*}
$$

Here $\gamma$ is a hyperparameter that controls the width of this Gaussianlike bumps. To gain an intuition for what the RBF kernel is doing, consider what prediction looks like in the perceptron:

$$
\begin{align*}
f(x) & =\sum_{n} \alpha_{n} K\left(x_{n}, x\right)+b  \tag{9.19}\\
& =\sum_{n} \alpha_{n} \exp \left[-\gamma\left\|x_{n}-z\right\|^{2}\right] \tag{9.20}
\end{align*}
$$

In this computation, each training example is getting to "vote" on the label of the test point $x$. The amount of "vote" that the $n$th training example gets is proportional to the negative exponential of the distance between the test point and itself. This is very much like an RBF neural network, in which there is a Gaussian "bump" at each training example, with variance $1 /(2 \gamma)$, and where the $\alpha_{n}$ s act as the weights connecting these RBF bumps to the output.

Showing that this kernel is positive definite is a bit of an exercise in analysis (particularly, integration by parts), but otherwise not difficult. Again, the proof is provided in the appendix.

So far, you have seen two bsaic classes of kernels: polynomial kernels $\left(K(x, v z)=(1+x \cdot z)^{d}\right)$, which includes the linear kernel $(K(x, z)=x \cdot z)$ and RBF kernels $\left(K(x, z)=\exp \left[-\gamma\|x-z\|^{2}\right]\right)$. The former have a direct connection to feature expansion; the latter to RBF networks. You also know how to combine kernels to get new kernels by addition. In fact, you can do more than that: the product of two kernels is also a kernel.

As far as a "library of kernels" goes, there are many. Polynomial and RBF are by far the most popular. A commonly used, but technically invalid kernel, is the hyperbolic-tangent kernel, which mimics the behavior of a two-layer neural network. It is defined as:

$$
\begin{equation*}
K^{(\tanh )}=\tanh (1+x \cdot z) \quad \text { Warning: not psd } \tag{9.21}
\end{equation*}
$$

A final example, which is not very common, but is nonetheless interesting, is the all-subsets kernel. Suppose that your $D$ features are all binary: all take values 0 or 1 . Let $A \subseteq\{1,2, \ldots D\}$ be a subset of features, and let $f_{A}(x)=\bigwedge_{d \in A} x_{d}$ be the conjunction of all the features in $A$. Let $\phi(x)$ be a feature vector over all such $A s$, so that there are $2^{D}$ features in the vector $\phi$. You can compute the kernel associated with this feature mapping as:

$$
\begin{equation*}
K^{\text {(subs) }}(x, z)=\prod_{d}\left(1+x_{d} z_{d}\right) \tag{9.22}
\end{equation*}
$$

Verifying the relationship between this kernel and the all-subsets feature mapping is left as an exercise (but closely resembles the expansion for the quadratic kernel).

### 9.5 Support Vector Machines

Kernelization predated support vector machines, but SVMs are definitely the model that popularized the idea. Recall the definition of the soft-margin SVM from Chapter 6.7 and in particular the optimization problem (6.36), which attempts to balance a large margin (small $\|\boldsymbol{w}\|^{2}$ ) with a small loss (small $\xi_{n} s$, where $\xi_{n}$ is the slack on the $n$th training example). This problem is repeated below:

$$
\begin{align*}
\min _{\boldsymbol{w}, b, \boldsymbol{\xi}} & \frac{1}{2}\|\boldsymbol{w}\|^{2}+C \sum_{n} \xi_{n}  \tag{9.23}\\
\text { subj. to } & y_{n}\left(\boldsymbol{w} \cdot \boldsymbol{x}_{n}+b\right) \geq 1-\xi_{n} \\
& \xi_{n} \geq 0
\end{align*}
$$

Previously, you optimized this by explicitly computing the slack variables $\xi_{n}$, given a solution to the decision boundary, $w$ and $b$. However, you are now an expert with using Lagrange multipliers
to optimize constrained problems! The overall goal is going to be to rewrite the SVM optimization problem in a way that it no longer explicitly depends on the weights $w$ and only depends on the examples $x_{n}$ through kernel products.

There are $2 N$ constraints in this optimization, one for each slack constraint and one for the requirement that the slacks are nonnegative. Unlike the last time, these constraints are now inequalities, which require a slightly different solution. First, you rewrite all the inequalities so that they read as something $\geq 0$ and then add corresponding Lagrange multipliers. The main difference is that the Lagrange multipliers are now constrained to be non-negative, and their sign in the augmented objective function matters.

The second set of constraints is already in the proper form; the first set can be rewritten as $y_{n}\left(\boldsymbol{w} \cdot \boldsymbol{x}_{n}+b\right)-1+\xi_{n} \geq 0$. You're now ready to construct the Lagrangian, using multipliers $\alpha_{n}$ for the first set of constraints and $\beta_{n}$ for the second set.

$$
\begin{align*}
\mathcal{L}(\boldsymbol{w}, b, \boldsymbol{\xi}, \boldsymbol{\alpha}, \boldsymbol{\beta})= & \frac{1}{2}\|\boldsymbol{w}\|^{2}+C \sum_{n} \xi_{n}-\sum_{n} \beta_{n} \xi_{n}  \tag{9.24}\\
& -\sum_{n} \alpha_{n}\left[y_{n}\left(\boldsymbol{w} \cdot \boldsymbol{x}_{n}+b\right)-1+\xi_{n}\right] \tag{9.25}
\end{align*}
$$

The new optimization problem is:

$$
\begin{equation*}
\min _{\boldsymbol{w}, b, \boldsymbol{\xi}} \max _{\boldsymbol{\alpha} \geq 0} \max _{\boldsymbol{\beta} \geq 0} \mathcal{L}(\boldsymbol{w}, \boldsymbol{b}, \boldsymbol{\xi}, \boldsymbol{\alpha}, \boldsymbol{\beta}) \tag{9.26}
\end{equation*}
$$

The intuition is exactly the same as before. If you are able to find a solution that satisfies the constraints (e.g., the purple term is properly non-negative), then the $\beta_{n}$ s cannot do anything to "hurt" the solution. On the other hand, if the purple term is negative, then the corresponding $\beta_{n}$ can go to $+\infty$, breaking the solution.

You can solve this problem by taking gradients. This is a bit tedious, but and important step to realize how everything fits together. Since your goal is to remove the dependence on $w$, the first step is to take a gradient with respect to $w$, set it equal to zero, and solve for $w$ in terms of the other variables.

$$
\begin{equation*}
\nabla_{w} \mathcal{L}=w-\sum_{n} \alpha_{n} y_{n} x_{n}=0 \quad \Longleftrightarrow \quad w=\sum_{n} \alpha_{n} y_{n} x_{n} \tag{9.27}
\end{equation*}
$$

At this point, you should immediately recognize a similarity to the kernelized perceptron: the optimal weight vector takes exactly the same form in both algorithms.

You can now take this new expression for $w$ and plug it back in to the expression for $\mathcal{L}$, thus removing $w$ from consideration. To avoid subscript overloading, you should replace the $n$ in the expression for
$w$ with, say, $m$. This yields:

$$
\begin{align*}
\mathcal{L}(b, \xi, \boldsymbol{\alpha}, \boldsymbol{\beta})=\frac{1}{2} \| & \sum_{m} \alpha_{m} y_{m} x_{m} \|^{2}+C \sum_{n} \xi_{n}-\sum_{n} \beta_{n} \xi_{n}  \tag{9.28}\\
& -\sum_{n} \alpha_{n}\left[y_{n}\left(\left[\sum_{m} \alpha_{m} y_{m} x_{m}\right] \cdot x_{n}+b\right)-1+\xi_{n}\right] \tag{9.29}
\end{align*}
$$

At this point, it's convenient to rewrite these terms; be sure you understand where the following comes from:

$$
\begin{align*}
\mathcal{L}(b, \xi, \boldsymbol{\alpha}, \boldsymbol{\beta})= & \frac{1}{2} \sum_{n} \sum_{m} \alpha_{n} \alpha_{m} y_{n} y_{m} x_{n} \cdot x_{m}+\sum_{n}\left(C-\beta_{n}\right) \xi_{n}  \tag{9.30}\\
& -\sum_{n} \sum_{m} \alpha_{n} \alpha_{m} y_{n} y_{m} x_{n} \cdot x_{m}-\sum_{n} \alpha_{n}\left(y_{n} b-1+\xi_{n}\right)  \tag{9.31}\\
= & -\frac{1}{2} \sum_{n} \sum_{m} \alpha_{n} \alpha_{m} y_{n} y_{m} x_{n} \cdot x_{m}+\sum_{n}\left(C-\beta_{n}\right) \xi_{n}  \tag{9.32}\\
& -b \sum_{n} \alpha_{n} y_{n}-\sum_{n} \alpha_{n}\left(\xi_{n}-1\right) \tag{9.33}
\end{align*}
$$

Things are starting to look good: you've successfully removed the dependence on $w$, and everything is now written in terms of dot products between input vectors! This might still be a difficult problem to solve, so you need to continue and attempt to remove the remaining variables $b$ and $\xi$.

The derivative with respect to $b$ is:

$$
\begin{equation*}
\frac{\partial \mathcal{L}}{\partial b}=-\sum_{n} \alpha_{n} y_{n}=0 \tag{9.34}
\end{equation*}
$$

This doesn't allow you to substitute $b$ with something (as you did with $\boldsymbol{w})$, but it does mean that the fourth term $\left(b \sum_{n} \alpha_{n} y_{n}\right)$ goes to zero at the optimum.

The last of the original variables is $\xi_{n}$; the derivatives in this case look like:

$$
\begin{equation*}
\frac{\partial \mathcal{L}}{\partial \xi_{n}}=C-\beta_{n}-\alpha_{n} \quad \Longleftrightarrow \quad C-\beta_{n}=\alpha_{n} \tag{9.35}
\end{equation*}
$$

Again, this doesn't allow you to substitute, but it does mean that you can rewrite the second term, which as $\sum_{n}\left(C-\beta_{n}\right) \xi_{n}$ as $\sum_{n} \alpha_{n} \xi_{n}$. This then cancels with (most of) the final term. However, you need to be careful to remember something. When we optimize, both $\alpha_{n}$ and $\beta_{n}$ are constrained to be non-negative. What this means is that since we are dropping $\beta$ from the optimization, we need to ensure that $\alpha_{n} \leq C$, otherwise the corresponding $\beta$ will need to be negative, which is not
allowed. You finally wind up with the following, where $x_{n} \cdot x_{m}$ has been replaced by $K\left(x_{n}, x_{m}\right)$ :

$$
\begin{equation*}
\mathcal{L}(\boldsymbol{\alpha})=\sum_{n} \alpha_{n}-\frac{1}{2} \sum_{n} \sum_{m} \alpha_{n} \alpha_{m} y_{n} y_{m} K\left(\boldsymbol{x}_{n}, \boldsymbol{x}_{m}\right) \tag{9.36}
\end{equation*}
$$

If you are comfortable with matrix notation, this has a very compact form. Let 1 denote the $N$-dimensional vector of all 1s, let $y$ denote the vector of labels and let $\mathbf{G}$ be the $N \times N$ matrix, where $\mathbf{G}_{n, m}=$ $y_{n} y_{m} K\left(x_{n}, x_{m}\right)$, then this has the following form:

$$
\begin{equation*}
\mathcal{L}(\alpha)=\alpha^{\top} \mathbf{1}-\frac{1}{2} \alpha^{\top} \mathbf{G} \alpha \tag{9.37}
\end{equation*}
$$

The resulting optimization problem is to maximize $\mathcal{L}(\boldsymbol{\alpha})$ as a function of $\alpha$, subject to the constraint that the $\alpha_{n}$ s are all non-negative and less than $C$ (because of the constraint added when removing the $\beta$ variables). Thus, your problem is:

$$
\begin{equation*}
\min _{\alpha}-\mathcal{L}(\boldsymbol{\alpha})=\frac{1}{2} \sum_{n} \sum_{m} \alpha_{n} \alpha_{m} y_{n} y_{m} K\left(x_{n}, x_{m}\right)-\sum_{n} \alpha_{n} \tag{9.38}
\end{equation*}
$$

subj. to $0 \leq \alpha_{n} \leq C$
One way to solve this problem is gradient descent on $\alpha$. The only complication is making sure that the $\alpha$ s satisfy the constraints. In this case, you can use a projected gradient algorithm: after each gradient update, you adjust your parameters to satisfy the constraints by projecting them into the feasible region. In this case, the projection is trivial: if, after a gradient step, any $\alpha_{n}<0$, simply set it to 0 ; if any $\alpha_{n}>C$, set it to $C$.

### 9.6 Understanding Support Vector Machines

The prior discussion involved quite a bit of math to derive a representation of the support vector machine in terms of the Lagrange variables. This mapping is actually sufficiently standard that everything in it has a name. The original problem variables $(w, b, \xi)$ are called the primal variables; the Lagrange variables are called the dual variables. The optimization problem that results after removing all of the primal variables is called the dual problem.

A succinct way of saying what you've done is: you found that after converting the SVM into its dual, it is possible to kernelize.

To understand SVMs, a first step is to peek into the dual formulation, Eq (9.38). The objective has two terms: the first depends on the data, and the second depends only on the dual variables. The first thing to notice is that, because of the second term, the $\alpha \mathrm{s}$ "want" to
get as large as possible. The constraint ensures that they cannot exceed $C$, which means that the general tendency is for the $\alpha$ s to grow as close to $C$ as possible.

To further understand the dual optimization problem, it is useful to think of the kernel as being a measure of similarity between two data points. This analogy is most clear in the case of RBF kernels, but even in the case of linear kernels, if your examples all have unit norm, then their dot product is still a measure of similarity. Since you can write the prediction function as $f(\hat{x})=\operatorname{sign}\left(\sum_{n} \alpha_{n} y_{n} K\left(x_{n}, \hat{x}\right)\right)$, it is natural to think of $\alpha_{n}$ as the "importance" of training example $n$, where $\alpha_{n}=0$ means that it is not used at all at test time.

Consider two data points that have the same label; namely, $y_{n}=$ $y_{m}$. This means that $y_{n} y_{m}=+1$ and the objective function has a term that looks like $\alpha_{n} \alpha_{m} K\left(x_{n}, x_{m}\right)$. Since the goal is to make this term small, then one of two things has to happen: either $K$ has to be small, or $\alpha_{n} \alpha_{m}$ has to be small. If $K$ is already small, then this doesn't affect the setting of the corresponding $\alpha \mathrm{s}$. But if $K$ is large, then this strongly encourages at least one of $\alpha_{n}$ or $\alpha_{m}$ to go to zero. So if you have two data points that are very similar and have the same label, at least one of the corresponding $\alpha$ s will be small. This makes intuitive sense: if you have two data points that are basically the same (both in the $x$ and $y$ sense) then you only need to "keep" one of them around.

Suppose that you have two data points with different labels:
$y_{n} y_{m}=-1$. Again, if $K\left(x_{n}, x_{m}\right)$ is small, nothing happens. But if it is large, then the corresponding $\alpha$ s are encouraged to be as large as possible. In other words, if you have two similar examples with different labels, you are strongly encouraged to keep the corresponding $\alpha$ s as large as $C$.

An alternative way of understanding the SVM dual problem is geometrically. Remember that the whole point of introducing the variable $\alpha_{n}$ was to ensure that the $n$th training example was correctly classified, modulo slack. More formally, the goal of $\alpha_{n}$ is to ensure that $y_{n}\left(\boldsymbol{w} \cdot \boldsymbol{x}_{n}+b\right)-1+\xi_{n} \geq 0$. Suppose that this constraint it not satisfied. There is an important result in optimization theory, called the Karush-Kuhn-Tucker conditions (or KKT conditions, for short) that states that at the optimum, the product of the Lagrange multiplier for a constraint, and the value of that constraint, will equal zero. In this case, this says that at the optimum, you have:

$$
\begin{equation*}
\alpha_{n}\left[y_{n}\left(w \cdot x_{n}+b\right)-1+\xi_{n}\right]=0 \tag{9.39}
\end{equation*}
$$

In order for this to be true, it means that (at least) one of the following must be true:

$$
\begin{equation*}
\alpha_{n}=0 \quad \text { or } \quad y_{n}\left(\boldsymbol{w} \cdot \boldsymbol{x}_{n}+b\right)-1+\xi_{n}=0 \tag{9.40}
\end{equation*}
$$

A reasonable question to ask is: under what circumstances will $\alpha_{n}$ be non-zero? From the KKT conditions, you can discern that $\alpha_{n}$ can be non-zero only when the constraint holds exactly; namely, that $y_{n}\left(\boldsymbol{w} \cdot \boldsymbol{x}_{n}+b\right)-1+\xi_{n}=0$. When does that constraint hold exactly? It holds exactly only for those points precisely on the margin of the hyperplane.

In other words, the only training examples for which $\alpha_{n} \neq 0$ are those that lie precisely 1 unit away from the maximum margin decision boundary! (Or those that are "moved" there by the corresponding slack.) These points are called the support vectors because they "support" the decision boundary. In general, the number of support vectors is far smaller than the number of training examples, and therefore you naturally end up with a solution that only uses a subset of the training data.

From the first discussion, you know that the points that wind up being support vectors are exactly those that are "confusable" in the sense that you have to examples that are nearby, but have different labels. This is a completely in line with the previous discussion. If you have a decision boundary, it will pass between these "confusable" points, and therefore they will end up being part of the set of support vectors.

### 9.7 Exercises

## Exercise 9.1. TODO...

