CS 189: Introduction to Machine Learning - Discussion 7

1. Kernels

a) Given a data point, \mathbf{x} , in \mathbb{R}^n , and the feature mapping ϕ corresponding to a quadratic kernel $k(\mathbf{x}, \mathbf{y}) = (\mathbf{x}^T \mathbf{y} + c)^2$, what is the dimensionality of $\phi(\mathbf{x})$? What is the dimensionality of the feature mapping corresponding to a Gaussian/RBF kernel? How do we deal with this high-dimensional data?

Solution: The quadratic kernel is defined as: $k(\mathbf{x}, \mathbf{y}) = (\mathbf{x}^T \mathbf{y} + c)^2$. $k(\mathbf{x}, \mathbf{y}) = (\sum_{i=1}^n x_i y_i + c)^2 = \sum_{i=1}^n (x_i^2) (y_i^2) + \sum_{i=2}^n \sum_{j=1}^{i-1} (\sqrt{2} x_i x_j) (\sqrt{2} y_i y_j) + \sum_{i=1}^n (\sqrt{2} c x_i) (\sqrt{2} c y_i) + c^2$ $\Phi(\mathbf{z}) = \langle x_n^2, \dots, x_1^2, \sqrt{2} x_n x_{n-1}, \dots, \sqrt{2} x_n x_1, \sqrt{2} c x_n, \dots, \sqrt{2} c x_1, c \rangle$

So, the dimension is $1 + n + n + {n \choose 2} = (n+2)(n+1)/2 \approx \frac{n^2}{2}$.

The Gaussian kernel corresponds to an infinite dimensional feature mapping. Deriving the kernel expansion is a bit difficult, but if interested, it can be found on page 39 of http://arxiv.org/pdf/0904.3664v1.pdf. Although explicitly mapping our data into these high-dimensional spaces would be computationally infeasible, kernels allow us to implicitly perform this feature map while still remaining in the original space.

b) Why might we prefer to use kernels?

Solution: Allows computing inner products instead of feature mapping, doesn't require explicit feature map, lets you stay in original space, etc.

2. Feature Selection

Feature selection is the process of selecting a subset of a dataset's original features for use in learning tasks. Any algorithm for this task has two parts: a search technique for proposing new feature subsets, and an evaluation measure which scores the different feature subsets. Wrapper methods involve retraining a model with these different subsets. Filter methods don't interact with the model, but rather just the features; they model interactions between the features and eliminate highly correlated ones. Embedded methods embed feature selection directly into the loss function.

a) Why might we want to perform feature selection?

Solution: There are a few reasons. By reducing the number of features, we can make models simpler and more interpretible. Imagine we had a rich dataset on people, and were trying to predict obesity. The most predictive features are height and weight (BMI), so ideally we would reduce the model to a function of these two variables, rather than a function of all of the collected features. This both helps to prevent overfitting and allows the researcher to more easily interpret the result. It also reduces training time.

b) Describe the simplest wrapper method you can think of for performing feature selection. What is its runtime? Can we do better?

Solution: Cross-validate models trained on every possible combination of features. This exhaustively checks all possible combinations. The runtime, unfortunately, is $O(2^d \cdot t)$, where d is the number of features and t is the time to train a model.

- c) Now, let's try the embedded method approach to feature selection.
 - i) What kind of penalty term corresponds directly to feature selection?
 - ii) What issues arise when optimizing a function with this penalty term?
 - iii) Describe two vectors such that their ℓ -1 norms are equal, but their ℓ -2 norms are not.
 - iv) Assume a method exists for optimizing non-differentiable functions, as long as they are convex. How can we approximate the penalty term from part i, while still maintaining an "optimizable" loss function?

Solution:

- i. We would like to penalize the number of features considered. This is equivalent to penalizing the number of nonzero features in the weight vector: the ℓ -0 norm.
- ii. The ℓ -0 norm is neither convex nor differentiable.

iii.
$$V_0 = (-0.5, -0.5), V_1 = (1, 0).$$

 $||V_0||_1 = |-0.5| + |-0.5| = 1$
 $||V_1||_1 |= |1| + |0| = 1.$

¹These are known as proximal gradient methods:

$$||V_0||_2 = \sqrt{|-0.5|^2 + |-0.5|^2} = 0.71$$

 $||V_1||_2 = \sqrt{|1|^2 + |0|^2} = 1$

iv. Making use of the existing optimization method, we can use the ℓ -1 norm instead. This norm is piecewise linear, and convex. As you can see from the previous part, the ℓ -1 norm allows sparser solutions than the ℓ -2 norm, which penalizes high values in any one dimension.

Further Improvements Weicheng, a previous TA, actually worked on these types of problems at Google this summer, and made use of two other types of norms, the truncated ℓ -1 norm, and the Lorentzian norm. These are described in detail at the following link:

http://people.csail.mit.edu/torralba/courses/6.869/lectures/lecture19/lecture19.pdf

3. Leave One Out Cross-Validation²

K-Fold Cross Validation works pretty well, but imagine a case where our training data is really scarce. We could get even more use out of our data by training on all data points except for one, then attempt to classify the remaining point. This corresponds to setting K=n in K-Fold Cross Validation, and is called Leave-One-Out-Cross-Validation (LOOCV). Although this would seem to require n model fits, we'll derive a method which saves us this expense in the **linear regression** setting.

Let's write the formal definition of LOOCV:

 $LOOCV = \sum_{i=1}^{n} (y_i - \hat{y_i}^{-i})^2$, where $\hat{y_i}^{-i}$ is the estimator of y_i with the i-th data point held-out when fitting the model.

a) After fitting our regression model, we have $\hat{y} = X\hat{w}$. Express \hat{y} in terms of the actual labels, i.e. find H such that $\hat{y} = Hy$, where H is an $n \times n$ matrix ³.

Solution: The optimal solution to linear regression, \hat{w} , is $(X^TX)^{-1}X^Ty$. So, $\hat{y} = X\hat{w} = X(X^TX)^{-1}X^Ty = Hy$, and $H = X(X^TX)^{-1}X^T$.

²Problem Set-Up: Credit to Tom Mitchell and Andrew W. Moore, CMU

³This matrix is known as the hat matrix, and as a number of useful statistical properties: https://en.wikipedia.org/wiki/Hat_matrix

b) By definition, \hat{y}^{-i} minimizes $\sum_{j\neq i} (y_j - \hat{y}_j^{-i})^2$. Prove \hat{y}^{-i} minimizes the squared error for z where

$$z_j = \begin{cases} y_j, & j \neq i \\ \hat{y_i}^{-i}, & j = i \end{cases}$$

Solution: From the definition, $\hat{y}^{-i} = \operatorname{argmin} \sum_{j \neq i} (y_j - \hat{y}_j^{-i})^2$.

However, $z_i = \hat{y}_i^{-i}$. So, $\underset{j \neq i}{\operatorname{argmin}} \sum_{j \neq i} (y_j - \hat{y}_j^{-i})^2 = \underset{j}{\operatorname{argmin}} \sum_{j} (z_j - \hat{y}_j^{-i})^2 = \hat{y}_i^{-i}$.

c) Prove $\hat{y_i}^{-i} = \hat{y_i} - H_{ii}y_i + H_{ii}\hat{y_i}^{-i}$

Solution: From part a), $\hat{y_i} = \sum_j H_{ij} y_j$. From part b), $\hat{y_i}^{-i} = \sum_j H_{ij} z_j$. Then, $\hat{y_i} - \hat{y_i}^{-i} = \sum_j H_{ij} (y_j - z_j) = H_{ii} (y_i - \hat{y_i}^{-i})$. So, $\hat{y_i}^{-i} = \hat{y_i} - H_{ii} y_i + H_{ii} \hat{y_i}^{-i}$.

d) Using this result, show that $LOOCV = \sum_{i=1}^{n} \left(\frac{y_i - \hat{y}_i}{1 - H_{ii}}\right)^2$. With this result, we can implement LOOCV while only fitting one regression model

Solution: $LOOCV = \sum_{i=1}^{n} (y_i - \hat{y_i}^{-i})^2$. Let's manipulate the result from part c): $\hat{y_i}^{-i}(1 - H_{ii}) = \hat{y_i} - H_{ii}y_i$. Thus, $\hat{y_i}^{-i} = \frac{\hat{y_i} - H_{ii}y_i}{(1 - H_{ii})}$. $LOOCV = \sum_{i=1}^{n} (y_i - \frac{\hat{y_i} - H_{ii}y_i}{(1 - H_{ii})})^2 = \sum_{i=1}^{n} (\frac{y_i - \hat{y_i}}{1 - H_{ii}})^2$.

e) Give an example of a dataset where LOOCV will give a poor estimate of the error.

Solution: In a dataset where every point is duplicated, the LOOCV will give an estimate of zero error. Since we only hold out one point at each step, there is a high likelihood that duplicates (or, at least, very near duplicates) may be contained in the training set. This results in inaccurately low error rates. As the dataset grows larger, the probability of near duplicates increases, and so LOOCV is best used in the small data setting.