Active Learning for Text Classification

Aktivní učení pro klasifikaci textů

Masters's Degree Project

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Contents

Introduction

Chapter 1

Introduction to Decision Theory

Easily speaking, decision can be defined as a set of actions. We are going to work with text data and solve text classification problem. Thus, our main purpose is to introduce precise mathematical abstractions for an automatic and optimal decision making.

Considering binary classification problem. We are seeking to find such set of actions that will assign each input value to its class.

We would like to commence our formal definition with the data. Let $\mathbf{x} \in \mathcal{X} \subset \mathbb{R}^n$ and $\mathbf{y} \in \mathcal{Y} = \{(0,1)^T, (1,0)^T\}$, where $\mathbf{x}$ are feature vectors of size n and $\mathbf{y}$ are its labels assigned to the data from space $\mathcal{X}$. Each value from space $\mathcal{Y}$ can be represented as a one hot representation that consist from ones and zeros $\mathbf{y} \in \{(0,1)^T, (1,0)^T\}$. Therefore, first class is represented as $\mathbf{y} = (1,0)^T$ and second class is represented as $\mathbf{y} = (0,1)^T$. As a good example of previous definition, $\mathbf{x}$ can be a text document (represented in a mathematical form in order to meet a definition above) and $\mathbf{y}$ can be its label that will show if the document is relevant or not. As seen from an example, label and text are forming a tuple. In this work we are also considering our data as tuples of variables $(\mathbf{x}, \mathbf{y}) \in \mathcal{X} \times \mathcal{Y}$.

1.1   Decision Theory

In this section we would like to define solution of a classification problem as finding such set of actions (decision) that minimizes loss. By loss we can understand classification error, entropy, etc.. This problem is still set up vaguely because neither decision nor loss was properly defined.

Each classification problem can be summarized as finding the optimal boundary that will split the dataset with respect to labels. The boundary can be defined as an action that is done with respect to some conditions. Let $a \in \mathcal{A}$ is an action and $\mathcal{A}$ is an action space. As mentioned previously, we would like to make a decision (action) that will split the dataset in two sets and we want this classification to be as good as possible. Quality of the classification can be measured with specific metrics which will be introduced in further sections. In our case we would like to minimize losses (error of classification). As a result, tandem of an action and loss minimization metric lets us to understand how good the action is. Subsequent step is an introduction of a loss function L. Loss function is dependent on action $a \in \mathcal{A}$. Formal definition of the loss function is shown next.

According to [2] loss function L can be defines as

$$L = L(\theta, a), \tag{1.1}$$

where $\theta \in \Theta$ is parameters' vector where Θ is parameters space and $a \in \mathcal{A}$ is an action where $\mathcal{A}$ is an action space. In [2], it is also said that parameter θ can be understand as the true state of nature.

Definition above shows us loss with respect to one action and one true state of nature. Thus, we would like to define expected loss function.

Definition 1. [2]If $\pi^*(\theta)$ is believed probability distribution of θ at the time of decision making, the *Bayesian expected loss* of an action a is

$$\rho(\pi^*, a) = \mathbb{E}_{\pi^*}[L(\theta, a)], \tag{1.2}$$

$$= \int_{\Theta} L(\theta, a)\pi^* d\theta. \tag{1.3}$$

Basing on definition 1.2, the optimal action is the one that minimizes expected loss, which is defined as

$$a^* = \operatorname*{argmin}_{a \in \mathcal{A}} \mathbb{E}_{\pi^*}[L(\theta, a)]. \tag{1.4}$$

However, how can we connect optimal action a^* with given data? Basing on the data definitions from previous part, we can assume that $\mathcal{X} \times \mathcal{Y}$ is an infinite set and $(\mathbf{x}, \mathbf{y})$ are samples from this set. Considering $\mathbf{x} \in \mathcal{X}$ and $\mathbf{y} \in \mathcal{Y}$ are random variables. If both $\mathbf{x}$ and $\mathbf{y}$ are random variables then tuple $(\mathbf{x}, \mathbf{y})$ is a sample from joint probability density function $p(\mathbf{x}, \mathbf{y})$. With the usage of conditional probability rule $p(\mathbf{x}, \mathbf{y})$ can be written as

$$p(\mathbf{x}, \mathbf{y}) = p(\mathbf{y}|\mathbf{x})p(\mathbf{x}). \tag{1.5}$$

At this part we are able to show which meaning this optimal action a^* has. In equation 1.5 $p(\mathbf{x})$ is given and $p(\mathbf{y}|\mathbf{x})$ is an unknown pdf that we want to estimate. The optimal action will lead us to an estimate of $p(\mathbf{y}|\mathbf{x})$.

1.2 Decision Theory for Supervised Learning

Supervised learning requires validation set of the data. Validation set is splitted on training and testing sets. Thus, we have to extend data definition.

Considering the data $\tilde{\mathbf{X}} \subset \mathcal{X}$ and its labels $\tilde{\mathbf{Y}} \subset \mathcal{Y}$. We define cartesian product $\tilde{\mathbf{X}} \times \tilde{\mathbf{Y}}$ as a validation set.

1.2.1 Decision Theory and Support Vector Machine Algorithm

In this subsection we will continue construction of the decision theory on the example of Support Vector Machine (SVM) method. For simplicity lets consider linearly separable dataset. From the theoretical perspective SVM constructs hyperplane in high dimensional space that separates two classes. In this case our decision (action) is a hyperplane that will separate two classes. Equation of the hyperplane can be written as $f(\mathbf{x}, \mathbf{w}, b) = \mathbf{w}^T \mathbf{x} + b$ where $\mathbf{w} \in \mathbb{R}^n$ is a set of hyperplane parameters and $b \in \mathbb{R}$ is a bias . As a result, action space is represented as $(\mathbb{R}^n, \mathbb{R}) = \mathcal{A}$ and as a consequence tuple $(\mathbf{w}, b) \in \mathcal{A}$. From this knowledge we can define $\theta = (\mathbf{x}, \mathbf{y})$ where tuple $(\mathbf{x}, \mathbf{y}) \in \mathcal{X} \times \mathcal{Y}$ is parameters and $\Theta = \mathcal{X} \times \mathcal{Y}$ is a parameters' space. However, for now we are limited on $\Theta = \tilde{\mathbf{X}} \times \tilde{\mathbf{Y}}$. Considering loss function 1.1 that can be written as

$$L = L(\mathbf{x}, \mathbf{y}, \mathbf{w}, b). \tag{1.6}$$

Following task is to understand how good is our action (hyperplane estimation) with respect to the dataset. We can choose different types of loss functions such as cross entropy, hinge loss, etc.. The most basic approach for SVM method is hinge loss function [12] which is defined as

$$L(\mathbf{x}, \mathbf{y}, \mathbf{w}, b) = \max(0, 1 - y\hat{y}(\mathbf{x}, \mathbf{w}, b)) \tag{1.7}$$

where $\hat{y}(\mathbf{x}, \mathbf{w}, b) = \mathbf{w}^T\mathbf{x} + b$ and $y = \mathbf{y}_1$.

In terms of SVM method we want to find such hyperplane that will label input values as a first class if it is "above" the hyperplane and as a second class if it is "below" the hyperplane. At this point very important assumption will be introduced. In order to find an optimal hyperplane we assume that the data $\tilde{\mathbf{X}}$ and its labels $\tilde{\mathbf{Y}}$ fully describe spaces $\mathcal{X}$ and $\mathcal{Y}$. Moreover, as mentioned earlier, we consider $\mathbf{x} \in \mathcal{X}$ and $\mathbf{y} \in \mathcal{Y}$ are random variables with joint probability density function $p(\mathbf{x}, \mathbf{y})$. We will also assume that $\forall i \in \{1, .., N\}$, $(\mathbf{x}_i, \mathbf{y}_i) \in \tilde{\mathbf{X}} \times \tilde{\mathbf{Y}}$ are independent identically distributed. As a result, with the usage of those data, probability density function $p(\mathbf{x}, \mathbf{y})$ can be estimated as

$$p(\mathbf{x}, \mathbf{y}) = \frac{1}{N}\sum_{i=1}^{N}\delta(\mathbf{x} - \mathbf{x}_i, \mathbf{y} - \mathbf{y}_i) \tag{1.8}$$

where $\delta(\mathbf{x} - \mathbf{x}_i, \mathbf{y} - \mathbf{y}_i)$ is Dirac delta function which is centered in $(\mathbf{x}_i, \mathbf{y}_i)$.

Using (1.2) we can evaluate expected loss function for SVM as follows

$$\mathbb{E}_{\pi^*}L = \int_{\mathcal{X}\times\mathcal{Y}} L(\mathbf{x}, \mathbf{y}, \mathbf{w}, b)p(\mathbf{x}, \mathbf{y})d(\mathbf{x}, \mathbf{y}),$$

$$= \int_{\mathcal{X}\times\mathcal{Y}} \max(0, 1 - y_1\hat{y}(\mathbf{x}, \mathbf{w}, b))\frac{1}{N}\sum_{i=1}^{N}\delta(\mathbf{x} - \mathbf{x}_i, \mathbf{y} - \mathbf{y}_i)d(\mathbf{x}, \mathbf{y}),$$

$$= \frac{1}{N}\sum_{i=1}^{N}\max(0, 1 - y_{1,i}\hat{y}(\mathbf{x}_i, \mathbf{w}, b))$$

where $\hat{y}(\mathbf{x}_i, \mathbf{w}, b) = \mathbf{w}^T\mathbf{x}_i + b$ and $y_{1,i}$ is first component of $i - $ th vector $\mathbf{y}_i$ Expect loss function for SVM can be written as

$$\rho(\mathbf{x}_i, \mathbf{w}, b) = \frac{1}{N}\sum_{i=1}^{N}\max(0, 1 - y_i(\mathbf{w}^T\mathbf{x}_i + b)). \tag{1.9}$$

1.2.2 Decision Theory and Algorithm Based on Neural Network Function

1.2.2.1 Neural Network

In this work we are using encoding of the input sample $\mathbf{x} \in \mathcal{X}$ and then applying Feed Forward Neural Network in order to predict the output value $\mathbf{y} \in \mathcal{Y}$. There are plenty ways how to provide encoding of the data sample $\mathbf{x}$ but we are interested here only in Feed Forward Neural Network algorithm that assigns input value to a specific class.

First layer of NN is defined as

$$\mathbf{a}_1 = \mathbf{W}_1^T\mathbf{x} + \mathbf{b}_1 \tag{1.10}$$

where $\mathbf{W}_1$ is matrix of weights and $\mathbf{b}_1$ is a vector of bias values. First layer is called input layer.

Further layers of NN are formed as

$$\mathbf{a}_k = \mathbf{W}_k^T f(\mathbf{a}_{k-1}) + \mathbf{b}_k, \; k = \{2, ..., K - 1\}. \tag{1.11}$$

As seen from equation 1.11, neurons from each layer (except of input layer) take linear combination of the neurons from the previous layer. Function f is an activation function. Activation function is defined as a non-decreasing, continuous function. The most usable activation functions are sigmoid, relu, elu, and hyperbolic tangence functions [3].

Output values are computed with

$$\hat{\mathbf{y}} = \mathbf{W}_K^T f(\mathbf{a}_{K-1}) + \mathbf{b}_K. \tag{1.12}$$

In most cases sigmoid and softmax functions for output layer are used. Softmax function is used for classification problems.

1.2.2.2 Decision Theory

Decision Theory construction for the algorithm, based on a neural network function, is mostly the same as in 1.2.1. However in this case our decision is to find estimate $\hat{\mathbf{y}} = \hat{\mathbf{y}}(\mathbf{x}, \Omega)$ of the probability density function $p(\mathbf{y}|\mathbf{x})$ where $\mathbf{x}$ is the input data, $\Omega = \{\mathbf{W}_1, ...\mathbf{W}_K, \mathbf{b}_1..., \mathbf{b}_K\}$ is a set of all neural network function parameters and biases. Action space $\mathcal{A}$ will be parameters' and biases' space of $\hat{\mathbf{y}}$. Same as in 1.2.1 we can define $(\mathbf{x}, \mathbf{y})$ are parameters of the loss function and $\mathcal{X} \times \mathcal{Y}$ is a parameters' space. Another example of loss functions that we will use is cross entropy loss function, which is defined as

$$L(\mathbf{x}, \mathbf{y}, \Omega) = -y_1 \ln\left(\hat{y}_1(\mathbf{x}, \Omega)\right) - y_2 \ln\left((\hat{y}_2(\mathbf{x}, \Omega)\right), \tag{1.13}$$

where $\mathbf{y} = (y_1, y_2)^T$ **and** $\hat{\mathbf{y}} = (\hat{y}_1, \hat{y}_2)^T$**.** We consider $\mathbf{x} \in \mathcal{X}$ and $\mathbf{y} \in \mathcal{Y}$ are random variables with joint probability density function $p(\mathbf{x}, \mathbf{y})$. With the usage of the given dataset where $\forall i \in \{1, .., N\}$, $(\mathbf{x}_i, \mathbf{y}_i) \in \tilde{\mathbf{X}} \times \tilde{\mathbf{Y}}$ are independent identically distributed we can approximate $p(\mathbf{x}, \mathbf{y})$ as (1.8). Applying definition 1, expected loss for the algorithm based on a neural network function is evaluated as

$$
\begin{aligned}
\mathbb{E}_{\pi^*} L &= \int_{\mathcal{X} \times \mathcal{Y}} L(\mathbf{x}, y, \Omega) p(\mathbf{x}, \mathbf{y}) d(\mathbf{x}, \mathbf{y}), \\
&= -\int_{\mathcal{X} \times \mathcal{Y}} \left(y_1 \ln\left(\hat{y}_1(\mathbf{x}, \Omega)\right) + y_2 \ln\left((\hat{y}_2(\mathbf{x}, \Omega)\right)\right) \frac{1}{N} \sum_{i=1}^N \delta(\mathbf{x} - \mathbf{x}_i, \mathbf{y} - \mathbf{y}_i) d(\mathbf{x}, \mathbf{y}), \\
&= -\frac{1}{N} \sum_{i=1}^N \left(y_1 \ln\left(\hat{y}_1(\mathbf{x}, \Omega)\right) + y_2 \ln\left((\hat{y}_2(\mathbf{x}, \Omega)\right)\right),
\end{aligned}
\tag{1.14}
$$

where $\mathbf{y} = (y_1, y_2)^T$ **and** $\hat{\mathbf{y}} = (\hat{y}_1, \hat{y}_2)^T$**.**

1.2.2.3 Parameters Estimation

In further sections we are going to introduce more methods based on Neural Networks which will slightly differ between each other. Thus, we would like to cover more theory around parameters estimation. The very simple but efficient method is gradient descent. This method is denoted with the usage of equation

$$\hat{\Omega}_{n+1} = \hat{\Omega}_n - \eta \nabla L(\mathbf{X}, \mathbf{Y}, \hat{\Omega}_n) \tag{1.15}$$

where $\hat{\Omega}_n$ is the $n-$th iteration value of gradient descent of NN weights and its biases that converges to $\hat{\Omega}$. Value of $\nabla L(\mathbf{X}, \mathbf{Y}, \hat{\Omega}_n)$ is a gradient of loss function and η is a value that in terms of NN is defined as learning rate. As a result described method finds the way to the function minimum.

The usual gradient descent is extremely efficient method for complex $\hat{\mathbf{y}}(\mathbf{x}, \Omega)$ but what if loss function has some local minimums? In this case, algorithm will stop iterating when it will find local minimum. Due to the complex functions $\hat{\mathbf{y}}(\mathbf{x}, \Omega)$, loss function of its approximation will be non convex with big amount of local minimums and maximums. Solution of this problem is that algorithm must jump over this local minimum. In terms of probability theory this tool is represented as a random value from some distribution. In comparison to standard gradient descent the key difference is that only one piece of data (minibatch) from the training dataset is used to calculate the step. The data samples (minibatch) is picked randomly at each step.

1.2.3 Decision Theory and Naive Bayes Algorithm

Naive Bayes algorithm is a bit different to the algorithm based on Neural Networks and SVM. In the case of Naive Bayes we want to estimate $p(\mathbf{W}|\mathbf{x}, \mathbf{y})$ where $\mathbf{W} = (\mathbf{w}_1, \mathbf{w}_2, ..., \mathbf{w}_n) \subset \mathcal{W}$ is an action ($a = \mathbf{W}$, $a \in \mathcal{A}$). The reason why we look for an estimate of the $p(\mathbf{W}|\mathbf{x}, \mathbf{y})$ but not $p(\mathbf{y}|\mathbf{x}, \mathbf{w})$ is due to the fact that in case of estimating $p(\mathbf{y}|\mathbf{x}, \mathbf{W})$ we would have to work with normalization constant would be dependent on the set of parameters $\mathbf{W}$. That fact would make our computations very complicated. Before continuing with a loss function construction we would like to go trough Naive Bayes (NB) method.

1.2.3.1 Naive Bayes

Assuming binary classification problem. With the usage of Bayes rule we can rewrite $p(\mathbf{W}|\mathbf{x}, \mathbf{y})$ as follows

$$p(\mathbf{W}|\mathbf{x}, \mathbf{y}) = \frac{p(\mathbf{y})p(\mathbf{x}|\mathbf{y}, \mathbf{W})p(\mathbf{W})}{\int_{\mathcal{W}} p(\mathbf{x}, \mathbf{y}|\mathbf{W})p(\mathbf{W})d\mathbf{W}} \tag{1.16}$$

where $\mathcal{W}$ is assumed as a space of $\mathbf{W}$.

Naive Bayes method introduces very strong assumption in equation 1.16. This assumption says that features of vector $\mathbf{x} = (x_1, x_2, ..., x_n)^T$ are conditionally independent. As a result estimation of $p(\mathbf{W}|\mathbf{x}, \mathbf{y})$ can be estimated as

$$\tilde{p}(\mathbf{W}|\mathbf{x}, \mathbf{y}) = \frac{1}{Z} p(\mathbf{y})p(\mathbf{W}) \prod_{i=1}^{n} (p(x_i|y_1, \mathbf{w}_i)^{y_1} p(x_i|y_2, \mathbf{w}_i)^{y_2}), \tag{1.17}$$

where $\mathbf{y} = (y_1, y_2)$, $\mathbf{W} = (\mathbf{w}_1, \mathbf{w}_2, ..., \mathbf{w}_n)$, Z is a normalizing constant.

1.2.3.2 Decision Theory

We want to maximize probability $\tilde{p}(\mathbf{W}|\mathbf{x}, \mathbf{y})$. As a result, using 1.17, loss function L will be represented as

$$L(\mathbf{y}, \mathbf{x}, \mathbf{w}) = -\log(\tilde{p}(\mathbf{W}|\mathbf{x}, \mathbf{y}), \tag{1.18}$$

$$= \log(Z) - \log(p(\mathbf{y})) - \log(p(\mathbf{W})) - \sum_{i=1}^{n} \log(p(x_i|y_1, \mathbf{w}_i)^{y_1} p(x_i|y_2, \mathbf{w}_i)^{y_2}). \tag{1.19}$$

Same as in 1.2.1 and 1.2.2 we will assume that we can approximate $p(\mathbf{x}, \mathbf{y})$ as 1.8. From this moment everything is ready for deriving expected loss function. Expected loss function for Naive Bayes method is derived as

$$\mathbb{E}_{\pi^*} L = \int_{X \times \mathcal{Y}} L(\mathbf{x}, \mathbf{y}, \mathbf{w}) p(\mathbf{x}, \mathbf{y}) d(\mathbf{x}, \mathbf{y}),$$

$$\int_{X \times \mathcal{Y}} \left(\xi(\mathbf{W}, \mathbf{y}) - \sum_{i=1}^{n} \log(p(x_i|y_1, \mathbf{w}_i)^{y_1} p(x_i|y_2, \mathbf{w}_i)^{y_2}) \frac{1}{N} \sum_{j=1}^{N} \delta(\mathbf{x} - \mathbf{x}_j, \mathbf{y} - \mathbf{y}_j) d(\mathbf{x}, \mathbf{y}),$$

$$\frac{1}{N} \sum_{j=1}^{N} \left(\xi_j(\mathbf{W}, \mathbf{y}_j) - \sum_{i=1}^{n} \log(p(x_{i,j}|y_{1,j}, \mathbf{w}_i)^{y_{1,j}} p(x_{i,j}|y_{2,j}, \mathbf{w}_i)^{y_{2,j}}) \right) \tag{1.20}$$

where $\xi(\mathbf{W}, \mathbf{y}) = \log(Z) - \log(p(\mathbf{y})) - \log(p(\mathbf{W}))$ and $\xi_j(\mathbf{W}, \mathbf{y}_j) = \log(Z) - \log(p(\mathbf{y}_j) - \log(p(\mathbf{W})).$

1.2.4 Decision Theory and Random Forest Algorithm

In order to work with random forests we must precisely define decision trees and only then construct theory for random forests.

1.2.4.1 Decision Tree

In this section we expect our decision three to give us an estimate $\hat{\mathbf{y}}(\mathbf{x}, \mathbf{w}) \in \{(0, 1)^T, (1, 0)^T\}$ where $\mathbf{w}$ is a vector that describes tree (depth, branches, etc.), $\mathbf{x} \in \mathcal{X}$ and $\mathbf{y} \in \mathcal{Y}$. It is very important to mention that for different trees $\mathbf{w}$ can have different dimensionality. Thus, for consistency we will assume that for all $\mathbf{w} \in \mathcal{W}$ exists upper bound, where $\mathcal{W}$ is redefined as a space of tree parameters. As a result we will make all $\mathbf{w}$ same length. If $\mathbf{w}$ has spare elements, they will be filled with zeros. Parameter space will be same as in 1.2.1-1.2.3, whereas action $a \in \mathcal{A}$ will be represented as $\mathcal{A} = \mathcal{W}$. In order to understand when our tree is optimal we can use zero-one loss function. Zero-one loss function is defined as

$$L(\mathbf{x}, \mathbf{y}, \mathbf{w}) = \begin{cases} 1, & \mathbf{y} \neq \hat{\mathbf{y}}(\mathbf{x}, \mathbf{w}) \\ 0, & \mathbf{y} = \hat{\mathbf{y}}(\mathbf{x}, \mathbf{w}) \end{cases}. \tag{1.21}$$

With the usage of the given data where $\forall i \in \{1, .., N\}$, $(\mathbf{x}_i, \mathbf{y}_i) \in \tilde{\mathbf{X}} \times \tilde{\mathbf{Y}}$ are independent identically distributed we can approximate $p(\mathbf{x}, \mathbf{y})$ as (1.8). As a result, expected loss function for decision tree can be derived as

$$\mathbb{E}_{\pi^*} L = \int_{\mathcal{X} \times \mathcal{Y}} L(\mathbf{x}, \mathbf{y}, \mathbf{w}) p(\mathbf{x}, \mathbf{y}) d(\mathbf{x}, \mathbf{y}),$$

$$= \int_{\mathcal{X} \times \mathcal{Y}} L(\mathbf{x}, \mathbf{y}, \mathbf{w}) \frac{1}{N} \sum_{i=1}^{N} \delta(\mathbf{x} - \mathbf{x}_i, \mathbf{y} - \mathbf{y}_i) d(\mathbf{x}, \mathbf{y}),$$

$$= \frac{1}{N} \sum_{i=1}^{N} L(\mathbf{x}_i, \mathbf{y}_i, \mathbf{w})$$

where $\sum_{i=1}^{N} L(\mathbf{x}_i, \mathbf{y}_i, \mathbf{w})$ is (1.21).

If we want to minimize expected loss we have to follow next steps. While construction decision tree we choose such feature $x_i \in (x_1, ..., x_n)^T = \mathbf{x}$ that will bring as the highest information about the system. This feature will form first layer, then we add another feature with the highest informational gain and construct second layer. Basing of this method we construct nodes and add more and more layers (branches).

In the following part we are going to work with a set of decision trees. For this purposes we will define our decision tree as $\hat{\mathbf{y}} = (T_1(\mathbf{x}, \mathbf{w}_l), T_2(\mathbf{x}, \mathbf{w}_l))^T$ where index l represents set of parameters for $l-$th three and indices $1, 2$ represent first and second value of the one-hot vector.

1.2.4.2 Random Forest

Considering $\mathbf{x} \in \mathcal{X}$ and $\mathbf{y} \in \mathcal{Y}$ as random variables with joint probability density function $p(\mathbf{x}, \mathbf{y})$. We will also assume that $\forall i \in \{1, .., N\}$, $(\mathbf{x}_i, \mathbf{y}_i) \in \tilde{\mathbf{X}} \times \tilde{\mathbf{Y}}$ are independent identically distributed.

With the usage of $\tilde{\mathbf{X}}$ and $\tilde{\mathbf{Y}}$ we will make $\{1, ..., L\}$, $L \in \mathbb{N}$ sets where $\forall l \in L$, $\tilde{\mathbf{X}}_l \subset \tilde{\mathbf{X}}$ and $\tilde{\mathbf{Y}}_l \subset \tilde{\mathbf{Y}}$. The data $\tilde{\mathbf{X}}_l$ and $\tilde{\mathbf{Y}}_l$ are created with random uniform sampling from $\tilde{\mathbf{X}}$ and $\tilde{\mathbf{Y}}$. We also want each subset to contain strictly 60% of the data from $\tilde{\mathbf{X}}$ and $\tilde{\mathbf{Y}}$. As a result parameter space for random forests will

form tuples of sets $(\tilde{\mathbf{X}}_l, \tilde{\mathbf{Y}}_l)$. Basing on this theory we will construct L decision trees $\hat{y} = T(\mathbf{x}, \mathbf{w}_l)$, where $\mathbf{x} \in \tilde{\mathbf{X}}_l$. As a result for $l - $ th decision tree

$$\mathbb{E}_{\pi^*} L = \frac{1}{N_l} \sum_{i=1}^{N_l} L(\mathbf{x}_{i,l}, \mathbf{y}_{i,l}, \mathbf{w}_l) \delta(\mathbf{x} - \mathbf{x}_{i,l}, \mathbf{y} - \mathbf{y}_{i,l}) \tag{1.22}$$

where $(\mathbf{x}_{i,l}, \mathbf{y}_{i,l}) \in (\tilde{\mathbf{X}}_l, \tilde{\mathbf{Y}}_l)$ and N_l is a number of the data in $\tilde{\mathbf{X}}_l$ and $\tilde{\mathbf{Y}}_l$. If we assume $\mathbf{w}_l$ as a random variable then L decision trees form samples from probability density function $p(y|\mathbf{x}, \mathbf{w})$. In other words

$$p(\mathbf{y}|\mathbf{x}, \mathbf{w}_l) = T_1(\mathbf{x}, \mathbf{w}_l)^{y_1} T_2(\mathbf{x}, \mathbf{w}_l)^{y_2} \tag{1.23}$$

where label $\mathbf{y}$ is written as a one-hot representation $\mathbf{y} = (y_1, y_2)^T$. Thus, we can say that classification probability $p(\mathbf{y}|\mathbf{x})$ can be written as

$$p(\mathbf{y}|\mathbf{x}) = \int_{\mathbf{w} \in \mathcal{A}} p(\mathbf{y}|\mathbf{x}, \mathbf{w}) p(\mathbf{w}) d\mathbf{w}. \tag{1.24}$$

where $\mathcal{A}$ is an action space. With the usage of samples $\mathbf{w}_l$ we can approximate $p(\mathbf{y}|\mathbf{x})$ as

$$p(\mathbf{y}|\mathbf{x}) = \frac{1}{L} \sum_{l=1}^{L} T_1(\mathbf{x}, \mathbf{w}_l)^{y_1} T_2(\mathbf{x}, \mathbf{w}_l)^{y_2} \tag{1.25}$$

where each decision tree $(T_1(\mathbf{x}, \mathbf{w}_l), T_2(\mathbf{x}, \mathbf{w}_l))^T$ is constructed with the usage of (1.22).

Before continuing with further sections we define output of Random Forest as $\hat{\mathbf{y}} = \hat{\mathbf{y}}(\mathbf{x}, \mathbf{W})$, where $\mathbf{W} = (\mathbf{w}_1, ..., \mathbf{w}_L)$ is set of parameters of specific Random Forest algorithm. We define vector $\hat{\mathbf{y}}$ as

$$\hat{\mathbf{y}} = \frac{1}{L} \sum_{l=1}^{L} (T_1(\mathbf{x}, \mathbf{w}_l), T_2(\mathbf{x}, \mathbf{w}_l)). \tag{1.26}$$

1.3   Decision Theory for Active Learning

As mentioned in previous sections $\tilde{\mathbf{X}} \times \tilde{\mathbf{Y}}$ is a validation set. When we finish a model training, we may think that we need more training data. Thus, we can choose the data from $\mathbf{X} \subset \mathcal{X} \backslash \tilde{\mathbf{X}}$. However it is important to understand that we have no labels for the set $\mathbf{X}$. We can ask for a help from an annotator that can give us those labels. We assume that getting labels needs some time and is very expensive.

Active learning problem is defined as a sequence of Supervised learning problems. Specifically, we assume that labels $\mathbf{y} \in \tilde{\mathbf{Y}}$ are available only for $\mathbf{x} \in \tilde{\mathbf{X}}$. We have the possibility to select an unlabeled element from $\mathbf{X}$ and ask for its label. Since it is expensive, we aim to have such questions that will maximize scores as fast as possible. Formally, we denote $J_0 = \{j_{01}, j_{02}, \dots j_{0N}\} = \{1, \dots, N\}$ the initial set of N available labels. using only the labeled data, the supervised learning task is defined on sets $\mathbf{X}_0 = \tilde{\mathbf{X}} = \{\mathbf{x}_i\}_{i \in J_0}$ and $\mathbf{Y}_0 = \tilde{\mathbf{Y}} = \{\mathbf{y}_i\}_{i \in J_0}$. This set is sequentially extended with new labels gained from $\mathbf{X}$. We consider a sequence of U questions $u = \{1, \dots, U\}$, in each question, we select an index j_u and ask to obtain the label $\mathbf{y}_{j_u}$ for data record $\mathbf{x}_{j_u}$. The index set and the data sets are extended as follows

$$J_u = \{J_{u-1}, j_u\}, \qquad \mathbf{X}_u = \{\mathbf{X}_{u-1}, \mathbf{x}_{j_u}\}, \qquad \mathbf{Y}_u = \{\mathbf{Y}_{u-1}, \mathbf{y}_{j_u}\}.$$

The task of active learning is to optimize the selection of indices j_u to reach as good classification metrics with as low number of questions as possible. As a result we have to define expected loss for each question u that will be dependent on the action and parameter spaces. In this case we can define our

action space $\mathcal{A}_u$ as a space of the data indices with respect to the parameters space Θ_u for each question u. Parameters space is defined as a set of possible parameters from a specific model. It is very important to understand that we will need not only one set of parameters but parameters distribution. We need parameters distribution because we will integrate over the parameters space. As an example, if we talk about SVM method, then parameters space for active learning problem will be defined as a set of weights that form a hyperplane. If we talk about the algorithm that is based on a Neural Network function, then parameters space of the active learning problem will form weights from neurons. We wanted to highlight that parameters space will be different for each problem but the idea for each algorithm is same .

As a result our task can be written as

$$j_u^* = \underset{j \in J \setminus J_u}{\operatorname{argmin}} (\mathbb{E}_{\pi_u^*} L^*) \tag{1.27}$$

where $\mathbb{E}_{\pi_u^*} L^*$ is expected loss that is dependent on an action given question u, and J is space of all indices. Expected loss for the active learning problem is defined as

$$\mathbb{E}_{\pi_u^*} L^* = \int_{\Theta_u} L^*(a, \theta) \pi_u^* d\theta \tag{1.28}$$

where $a \in \mathcal{A}_u$ and $\theta \in \Theta_u$ and L^* is a loss function for the active learning problem. Character "$*$" is used only for distinguishing active learning loss from the loss function which is used for different models. We will specify action space because it will be the same for all models that are used in the active learning section. Action space $\mathcal{A}$ is a set of possible indices $j_u \in J_u$ where u is a specific question. Thus, (1.28) can be written as

$$\mathbb{E}_{\pi_u^*} L^* = \int_{\Theta_u} L^*(j_u, \theta) \pi_u^* d\theta. \tag{1.29}$$

Using this approach we will be able sequentially select indices from $\mathbf{X}$ and ask for a label from $\mathbf{Y}$ that will help us to get higher scores faster than in the case of random choice of indices.

1.3.1 Bayesian Approach of Classifiers' Parameters Sampling

Considering that $\mathbf{y} \in \mathcal{Y}$. Let $\hat{\mathbf{y}} = \hat{\mathbf{y}}(\mathbf{x}_{j_u}, \theta_u)$ is an estimate of $\mathbf{y}$. However, in this case output estimate $\hat{\mathbf{y}}$ is represented as a vector of probabilities that $\mathbf{x}_{j_u}$ is assigned to different classes. As an example for well trained binary classifier, for specific $\mathbf{x}$ that is assigned to $\mathbf{y} = (1, 0)^T$, classifiers estimate of $\mathbf{x}$ can be $\hat{\mathbf{y}} = (0.95, 0.05)^T$. It is very interesting that before we can solve the optimization problem with choosing the index j_u we have to solve the optimization problem of finding $\hat{\mathbf{y}}$. This leads us to supervised learning models that we have discussed in previous sections.

In this section we would like to construct theory around π_u^* from equation 1.29. Mentioned distribution is a distribution of the models' parameters given the training data that can be written as

$$\pi_u^* = p_u(\theta_u | \mathbf{X}_u, \mathbf{Y}_u). \tag{1.30}$$

We do not have explicit form of the pdf. However, we assume that we have samples Q_u samples $\theta_{u,q} \in \{1_u, ..., Q_u\}$ from $p_u(\theta_u | \mathbf{X}_u, \mathbf{Y}_u)$. As a result $p_u(\theta_u | \mathbf{X}_u, \mathbf{Y}_u)$ can be approximated as

$$\pi_u^* = \frac{1}{Q_u} \sum_{q=1}^{Q_u} \delta(\theta_u - \theta_u^{(q)}), \tag{1.31}$$

where $\delta(\theta_u - \theta_u^{(q)})$ is Dirac delta function centered in $\theta_u^{(q)}$.

1.3.1.1 Parameters Sampling Based on Training Data Subsets

This method is quite general and can be applied to all types of classifiers in this work (Random Forest, SVM, Neural Network). The idea is very simple. We consider that some data samples in training dataset $\tilde{\mathbf{X}} \times \tilde{\mathbf{Y}}$ are noise corrupted. Thus, its obvious that we do not want our models to learn from noise corrupted data. As a result, we would like to randomly sample Q_u subsets from $\tilde{\mathbf{X}}$ with their labels from $\tilde{\mathbf{Y}}$. Lets rewrite it in more mathematical form.

Assuming N_u is amount of samples in $\mathbf{X}_u$. Let $Z_u = \{z_1, ..., z_{N_u^{sub}}\} \subset J_u$, where $N_u^{sub} < N_u$. Let $p(\theta_u | \mathbf{X}_u, \mathbf{Y}_u, Z_u)$ is a probability of model parameters θ_u given $\mathbf{X}_u, \mathbf{Y}_u$ and Z_u. The conditioning in the pdf is defined as restriction of sets $\mathbf{X}_u$, $\mathbf{Y}_u$ on indices from Z_u. As a result we can approximate 1.30 as

$$\pi_u^* = p(\theta_u | \mathbf{X}_u, \mathbf{Y}_u) \tag{1.32}$$

$$= \int p(\theta_u | \mathbf{X}_u, \mathbf{Y}_u, Z_u) p(Z_u) dZ_u \tag{1.33}$$

$$= \frac{1}{Q_u} \sum_{q=1}^{Q_u} p(\theta_u | \mathbf{X}_u, \mathbf{Y}_u, Z_u^{(q)}) \tag{1.34}$$

$$= \frac{1}{Q_u} \sum_{q=1}^{Q_u} \delta(\theta_u - \theta_u^{(q)}). \tag{1.35}$$

Sampling from $p(Z_u)$ is very simple. The only thing that must be predefined is N_u^{sub}. After training the model using $\mathbf{X}_u$ and Y_u under restriction Z_u, vector of model parameters will represent a single sample from π_u^*.

1.3.1.2 SGLD

Unlike previous section method, SGLD sampling is designed only for neural networks based classifiers. SGLD modifies Neural Network learning algorithm by adding noise in Stochastic Gradient descent. In comparison to Normal Feed Forward Neural Network in [13] is introduced Stochastic Gradient Langevin Dynamics Neural Network (SGLD). This is type of Bayesian Neural Network algorithm. In comparison to normal Feed Forward NN, [13] proposes to add gaussian noise to the system while doing gradient descent. With the usage of Bayes rule. we can rewrite 1.30 as

$$\begin{aligned} \pi_u^* &= p(\theta_u | \mathbf{X}_u, \mathbf{Y}_u) \\ &\propto p(\mathbf{X}_u, \mathbf{Y}_u | \theta_u) p(\theta_u). \end{aligned} \tag{1.36}$$

Next, we can approximate 1.36 as

$$\pi_u^* = \frac{1}{Q_u} \sum_{q=1}^{Q_u} p(\mathbf{X}_u, \mathbf{Y}_u | \theta_u^{(q)}) \tag{1.37}$$

$$= \frac{1}{Q_u} \sum_{q=1}^{Q_u} \delta(\theta_u - \theta_u^{(q)}).$$

However, for this case θ must meet some constraints. If we want to sample parameters from its distribution θ must be independent identically distributed. This is shown and proved in [13].

1.3.1.3 DENFI

In case of DENFI algorithm we can approximate 1.30 with the usage of equations 1.36 and 1.37. The equations are same but sampling from $p(\theta_u)$ is different. In this work we are using modification of the initial DENFI algorithm. The idea of this algorithm is to train ensemble of Q_u Feed Forward Neural Networks with the usage of stochastic gradient descent. In theory, each neural network will find different local minimum due to different initial weights of neurons and stochastic gradient descent. The beauty of this algorithm is shown in further training iterations that will be described and shown in further sections. For now, we are only interested in parameters sampling from $p(\theta_u)$, that has been already covered and described.

1.3.2 Active Learning Loss Function

1.3.2.1 Entropy Based Active Learning Loss

First approach of defining Active Learning loss function is negative entropy. Basing of the formal definition of the entropy we can write it as

$$-H(\hat{\mathbf{y}}|\mathbf{x}_{j_u}, \theta_u) = \sum_{r=1}^{R} \hat{y}_r(\mathbf{x}_{j_u}, \theta_u) \log(\hat{y}_r(\mathbf{x}_{j_u}, \theta_u)), \tag{1.38}$$

where $\hat{y}_r$ is $r-$th element of the output estimate $\hat{\mathbf{y}}$ and θ is a vector of parameters for specific model. As done in Passive Learning sections we want to find expected loss based on entropy function.

With the usage of previous knowledge we can derive expected entropy loss as

$$\mathbb{E}_{\pi_u^*} L^* = \int_{\Theta_u} -H(\hat{\mathbf{y}}|\mathbf{x}_{j_u}, \theta_u) p_u(\theta_u|\mathbf{X}_u, \mathbf{Y}_u) d\theta_u$$

$$= \int_{\Theta_u} -H(\hat{\mathbf{y}}|\mathbf{x}_{j_u}, \theta_u) \frac{1}{Q_u} \sum_{q=1}^{Q_u} \delta(\theta_u - \theta_{u,q}) d\theta_u$$

$$= \frac{1}{Q_u} \sum_{q=1}^{Q_u} -H(\hat{\mathbf{y}}|\mathbf{x}_{j_u}, \theta_{u,q}))$$

$$= \frac{1}{Q_u} \sum_{q=1}^{Q_u} \sum_{r=1}^{R} \hat{y}_r(\mathbf{x}_{j_u}, \theta_{u,q}) \log(\hat{y}_r(\mathbf{x}_{j_u}, \theta_{u,q})). \tag{1.39}$$

As a result, minimization of given expected loss will lead us to a sample with the highest entropy.

1.3.2.2 False Positive Sampling Loss

TBD

1.3.3 Active Learning for Random Forests Algorithm

In Supervised Learning section 1.2.4 we have derived that $p(\mathbf{y}|\mathbf{x})$ for Random Forest (RF) algorithm is written as 1.25. Active Learning algorithm requires distribution over the parameters of the Random Forest algorithm. We will solve this problem the way that we will get samples from π_u^* and then approximate probability distribution as (1.30).

In order to estimate samples from π_u^* we define Ensemble Random Forest Algorithm. That means that we will use Q_u Random Forests in each step of Active Learning algorithm. In this case parameters

of each $\hat{\mathbf{y}}_{q_u}$, where $\hat{\mathbf{y}}$ is Random Forest one-hot represented output and $q_u \in \{1_u, ..., Q_u\}$, will be iid. Ofcourse RF algorithm is already ensemble of decision trees but in this case we create ensemble algorithm from ensemble algorithms. Previously we defined for each decision tree T_l where $l \in \{1, ..., L\}$, that $\mathbf{w}_l$ is a vector of parameters of $l-$th decision tree. Thus, let $\mathbf{W} = (\mathbf{w}_1, ..., \mathbf{w}_L)$ is set of parameters of specific Random Forest algorithm. Thus, for $\hat{\mathbf{y}}_{q_u}$ we have $\mathbf{W}_{q_u} = (\mathbf{w}_{q_u,1}, ..., \mathbf{w}_{q_u,L})$. As a result we can approximate π_u^* as

$$\pi_u^* = \frac{1}{Q_u} \sum_{q_u=1}^{Q_u} \delta(\mathbf{W} - \mathbf{W}_{q_u}). \tag{1.40}$$

Another important point is to define $\hat{\mathbf{y}}_{q_u} = \hat{\mathbf{y}}(\mathbf{x}_{j_u}, \mathbf{W}_{q_u})$. We formally define each element $\hat{y}_{r,q_u}$ of vector $\hat{\mathbf{y}}_{q_u}$ as

$$\hat{y}_{r,q_u} = \frac{1}{L} \sum_{l=1}^{L} T_r(\mathbf{x}_{j_u}, \mathbf{w}_{q_u,l}),$$

where r is $r-$th class T_r is $r-$th element of a decision tree vector.

With the usage of the theory from section 1.3.2 we can write expected loss for Active Learning Random Forest Algorithm as

$$\mathbb{E}_{\pi_u^*} L_u^* = \frac{1}{Q_u} \sum_{q=1}^{Q_u} \sum_{r=1}^{R} \hat{y}_r(\mathbf{x}_{j_u}, \mathbf{W}_{q_u}) \log (\hat{y}_r(\mathbf{x}_{j_u}, \mathbf{W}_{q_u})). \tag{1.41}$$

Chapter 2

Natural Language Processing Theory

2.1 Text Representation

According to [5], Natural Language Processing (NLP) is a theoretically motivated range of computational techniques for analyzing and representing naturally occurring texts at one or more levels of linguistic analysis for the purpose of achieving human-like language processing for a range of tasks or applications.

In this work we are focused on two techniques such as TF-IDF [11] and Fast Text Word Embeddings [7]. These methods are used for representation of text in a mathematical form (vectors, matrices). Even though TF-IDF is quite old method for text representation, it is still widely used. However, primary method, that used in the theses is Fast Text Word Embeddings. In this project we are working with text documents (articles and tweets) and their labels. In the beginning of chapter 1 we defined value $\mathbf{x} \in \mathcal{X}$ as text features vector. By features vector we mean any kind of text encoding (TF-IDF, Fast Text Word Embedding, etc..).

2.1.1 TF-IDF

Term Frequency - Inverse Document Frequency (TF-IDF) is extremely powerful tool. This text encoding tool is quite simple and powerful. Method's advantage is its popularity. Plenty of packages in different programming languages have implementations of this algorithm. As mentioned in the name of this method, it is composed from two parts Term Frequency and Inverse Document Frequency. Term Frequency is defined as

$$TF(t,d) = \frac{f_{t,d}}{\sum_{t'} f_{t',d}},$$

where $f_{t,d}$ is number of times of word t in a document d. Inverse Document Frequency is defined as

$$IDF(t,d) = \log \frac{|D|}{|\{d \in D : t \in d\}|},$$

where numerator stand for total number of documents in the corpus and denominator is number of documents where the term t appears. We are assuming using only words that from corpus D. Thus, the denominator is always greater than zero.

Finally,

$$TF - IDF(t,d) = TF(t,d) \cdot IDF(t,d).$$

N.B

2.1.1.1   TF-IDF and Information Theory

In this part is shown the connection of TF-IDF to Information theory [1]. Lets first take a look on documents' entropy given word t,

$$
\begin{aligned}
H(D|T = t) &= - \sum_d p(d|t) \log p(d|t) \\
&= \log \frac{1}{|\{d \in D : t \in D\}|} \\
&= - \log \frac{|\{d \in D : t \in D\}|}{D|} + \log |D| \\
&= -IDF(t, d) + \log |D|,
\end{aligned}
\tag{2.1}
$$

where D is a documents' random variable and T is words' random variable. Equation 2.1 is correct under condition that we have no duplicate documents in the text corpus. Next step is to derive an equation of mutual information of documents and words as follows

$$
\begin{aligned}
M(D, T) &= H(D) - H(D|T) \\
&= - \sum_d p(d) \log p(d) - \sum H(D|T = t) \cdot p(t)_t \\
&= \sum_t p(t) \cdot \left(\log \frac{1}{|D|} + IDF(t, d) - \log |D| \right) \\
&= \sum_t p(t) \cdot IDF(t, d) \\
&= \sum_{t,d} p(t|d) \cdot p(d) \cdot IDF(t, d) \\
&= \frac{1}{|D|} \sum_{t,d} TF(t, d) \cdot IDF(t, d).
\end{aligned}
\tag{2.2}
$$

As seen from 2.2 TF-IDF has really good explanatory definition based on Information Theory. As a result, it is one more advantage of this method usage. However, here is one big disadvantage that can be very crucial. The higher amount of words is, the bigger and sparser vectors, that represent each document, will be.

2.1.2   Fast Text and CBOW Word Embeddings

Term word embedding means a set of language modeling and feature learning techniques in natural language processing where words or phrases from the vocabulary are mapped to vectors of real numbers. Nowadays exist plenty of word embedding methods based on neural networks and co-occurrence matrices. Word embeddings are used as pretrained models. Words' encoding is used in order to encode the text and then text encoding is used for different purposes such as classification, clustering, etc..

The principle of word embeddings based on neural networks is explained in this section. We decided to describe Continuous Bag of Words Model (CBOW), because Fast Text word embeddings model is a modification of this method and CBOW covers all main theoretical aspects.

2.1.2.1   CBOW Word Embeddings

[6]

We would like to treat text {"The", "cat", 'over", "the', "puddle"} as a context and from these words, be able to predict or generate the center word "jumped". This type of model is a Continuous Bag of Words (CBOW) Model. Before we continue with more theoretical part, it is good to mention that mathematical notation defined here is only used for this section and has no common with the same names of the variables that were defined in the beginning of this theses. First, we want to set up our known parameters. Let the known parameters in our model be the sentence represented by one-hot word vectors. The input one hot vectors or context we will represent with an $\mathbf{x}^{(c)}$. And the output as $\mathbf{y}^{(c)}$ and in the CBOW model, since we only have one output, so we just call this $\mathbf{y}$ which is the one hot vector of the known center word. Now let's define our unknowns in our model. We create two matrices, $\mathcal{V} \in \mathbb{R}^{n \times |V|}$ and $\mathcal{U} \in \mathbb{R}^{|V| \times n}$. Where n is an arbitrary size which defines the size of our embedding space. $\mathcal{V}$ is the input word matrix such that the $i-$th column of $\mathcal{V}$ is the $n-$dimensional embedded vector for word w_i when it is an input to this model. We denote this $n \times 1$ vector as $\mathbf{v}_i$. Similarly, $\mathcal{U}$ is the output word matrix. The $j-$th row of $\mathcal{U}$ is an $n-$dimensional embedded vector for word w_j when it is an output of the model. We denote this row of $\mathcal{U}$ as $\mathbf{u}_j$.

For this method sequence of actions can be written as follows:

- We generate our one hot word vectors $(\mathbf{x}^{(c-m)}, ..., \mathbf{x}^{(c-1)}, \mathbf{x}^{(c+1)}, ..., \mathbf{x}^{(c+m)})$ for the input context of size $2m$.

- We get our embedded word vectors for the context $(\mathbf{v}_{c-m} = \mathcal{V}\mathbf{x}^{(c-m)}, \mathbf{v}_{c-m+1} = \mathcal{V}\mathbf{x}_{(c-m+1)}, ..., \mathbf{v}_{c+m} = \mathcal{V}\mathbf{x}_{(c+m)})$

- Average these vectors to get $\tilde{\mathbf{v}} = \frac{\mathbf{v}_{c-m} + \mathbf{v}_{c-m+1} + ... + \mathbf{v}_{c+m}}{2m}$

- Generate a score vector $\mathbf{z} = \mathcal{U}\tilde{\mathbf{v}}$

- Turn the scores into probabilities

$$\hat{\mathbf{y}} = \text{softmax}(\mathbf{z}) \tag{2.3}$$

- We desire our probabilities generated, $\hat{\mathbf{y}}$, to match the true probabilities, $\mathbf{y}$, which also happens to be the one hot vector of the actual word.

Described method can be interpreted as a feed forward neural network with one hidden layer that do not uses activation function. As a loss function for this algorithm can be chosen cross-entropy loss function

$$L = \sum_{i=1}^{|V|} \mathbf{y_i} \log(\hat{\mathbf{y}_i}) \tag{2.4}$$

where $\hat{\mathbf{y}}$ is sofmatx (2.3) function.

2.1.2.2 Fast Text Word Embeddings

As mentioned previously, Fast Text method is a CBOW modification. Main modification is that Fast Text is taking into account not only words but also suffixes of words. The words are splitted into suffixes and as a result they can handle understanding of the context better.

In this theses we used pretrained Fast Text models [7] consisting of 1 million word vectors trained on Wikipedia 2017, UMBC webbase corpus and statmt.org news dataset (16B tokens).

Chapter 3

Data and Results Validation Theory

3.1 Results Validation Theory

When the models are implemented and trained we have to compare them. This part is very important because we want to define such metrics that will not be biased and which will have high discriminability. In this project experiments are separated into two parts. First part is supervised classification with big amount of data. This is done for understating what is the maximal upper bound of specific classifiers. These upper bounds are used as maximums which our active learning algorithms must be converging to.

3.1.1 Receiving Operating Curve metric

In section 1 we mentioned that we are limiting our problem only on binary classification. Plenty of metrics such as recall, accuracy, precision, etc. exists for binary classification. However we decided to find a metic that is able to unify all metrics discussed before and do not underperform each of them. For these purposes we chose Receiving Operating Curve metric. ROC visualizes the tradeoff between true positive rate (TPR) and false positive rate (FPR). This means that for every threshold, we are able to calculate TPR and FPR and plot it on one figure.

We are working with balanced datasets. Thus, there is no problem in using ROC metrics. ROC metric is also very good when we care equally about positive and negative class. Another advantage is that if we notice small changes in ROC it will not result in big changes in other binary classification metrics.

Receiving Operating Curve lets us to calculate area under the curve (AUC). The probabilistic interpretation of ROC score is that if a positive case and a negative case are chosen randomly, the probability that the positive case outranks the negative case according to the classifier is given by the AUC.

3.1.2 Supervised Learning Results Validation

As mentioned above, supervised learning results are used as maximal upper bound of specific classifiers. In order to make results statistically valid we used k-fold cross validation. For each batch from k-fold cross validation we calculated both ROC and AUC. As an output result of a classifier performance we calculated mean value through all ROC and AUC results. All results are calculated with respect to balanced data classes.

3.1.3 Active Learning Results Validation

Active learning model evolution is based on supervised learning algorithms that are sequentially retrained. Thus, we are not able to display ROC for active learning algorithms because amount of results

is too big. We decided to aggregate results and display evolution of AUC metric for each step of active learning sequence. AUC sequences can be well compared between different classifiers. Another aspect of data validation is making the results statistically significant. We are not able to use k-fold cross validation for active learning algorithms. Therefore we active learning algorithm $H \in \mathbb{N}$ times. Due to the fact of random initializations, we are able to determine uncertainty bounds that are calculated as standard deviations to mean value.

3.2 Data

This chapter is dedicated dataset description. We used two datasets for algorithms trainings and testings. We consider these datasets big and diverse enough for getting unbiased results. We took into account size of texts (articles and tweets), diversity and at the same time similarity of topics.

3.2.1 HuffPost 200k Articles Dataset

HuffPost 200k Articles Dataset is publicly available at Kaggle competition webpage and can be found as News Category Dataset [8] here `https://www.kaggle.com/rmisra/news-category-dataset`. Described dataset contains around 200k news headlines from the year 2012 to 2018 obtained from Huff-Post. Following dataset includes url address and label to each article. HuffPost dataset has 200k articles assigned to 41 categories. We used only 10 categories. We are interested in binary classification, as a result we have to make pairs from chosen categories. These categories and pairs are listed in table 3.1.

Tuple Id	Category Pairs	
1	Crime	Good News
2	Sports	Comedy
3	Politics	Business
4	Science	Tech
5	Education	College

Table 3.1: HuffPost Dataset Categories which were chosen for algorithms' training and testing

Due to the fact that there is no raw article included in the dataset, we used url links in order to find the articles. For each category we scraped 500 original articles from `www.huffpost.com`. Listed categories are chosen with respect to diversity and classification complexity. We sorted the categories in table 3.1 with respect to ascending classification complexity order. By classification complexity we mean two sets intersections in feature spaces. If the classification complexity is high, then majority of feature space dimensions have intersections between two datasets. Thus, it is harder to find such set of features that can be used for high classification performance.

3.2.2 1600k Tweets Dataset

Another dataset that is used in this work is 1600k Tweets dataset which is publicly available and is also taken from Kaggle competition webpage. The dataset can be found as sentiment140 dataset [4] here `https://www.kaggle.com/kazanova/sentiment140`. This dataset contains 1,600,000 tweets extracted using the twitter api. The tweets have been annotated as negative (0), positive (4) and they are used for sentiment detection. Same as with HuffPost dataset we used 500 records for each category for training and testing purposes. The reason of choosing this dataset is that we wanted to show how our algorithms can handle data that consist of little texts.

Chapter 4

Project Implementation and Architecture

Even though this work is quite theoretical with experiments that proof theoretical concepts, we assume implementational part interesting as well. In this chapter we show the architecture of the project and explain how different dependencies cooperate with each other. The codebase of this project can be easily found at `https://github.com/sahanmar/Peony`. We expect this project going to continue grow and will be used not only in terms of master theses.

This project was written in Python 3.7 programming language with the usage of Conda environment. The project combines a lot of different tools and programs such as Docker, Docker-Compose, MongoDb, Jupyter, etc..

In this theses we used two main components that represent the database and computational part. We tried to unify all methods as much as possible and make the utilization process very easy.

4.1   Database

In order to make everything consistent and let the models work with the same input and output format we decided to create a database that will store all the data in JSON format. This unification let us connect the database to machine learning and visualizing components. In this project we decided to work with NoSQL database. Our choice was MongoDb. The reason why we have chosen MongoDb is because of its simplicity and possibility of maintaining through Docker. Since Docker and MongoDb is perfect combination, the database can be deployed with two lines of code through Docker-Compose as explained in documentation on GitHub. Of course it is easier to use MongoDb without Docker but our motivation was measured on simplicity of creating and working with the database. All experiments were run on Google Cloud Platform Virtual Machine instance. Thus, we could start working with the models right away without any complications with installation.

4.1.1   MongoDb Data Format

MongoDb represents the data in BSON format behind the scenes but we are send and get there JSON format data. Despite the fact that we are having different text datasets that we store in the database, we decided to create unified JSON scheme that will let us to preserve the structure of the data stored in MongoDb. JSON schema of how the data are stored and what a user will get as an output from a database is shown in figure 4.1. Deeper explanation of JSON schema can be found at `https://json-schema.org/understanding-json-schema/`.

```json
{
  "title": "Database",
  "type": "object",
  "properties": {
    "datasetName": {
      "type": "string",
      "description": "Name of the dataset"
    },
    "datasetId": {
      "type": "int",
      "description": "Unique hash id that will be created automatically"
    },
    "record": {
      "type": "object",
      "description": "All information about an instance",
      "properties": {
        "id": {
          "type": "string",
          "description": "Unique hash id that will be created automatically
            "
        },
        "snippet": {
            "type": "string",
            "description": "Snippet of a text. Can be empty"
        },
        "text": {
          "type": "object",
          "description": "Text instance that is used for a model",
          "properties" : {
            "title": {
              "type": "string",
              "description": "Title of a text. Can be empty"
            },
            "body": {
              "type": "string",
              "description": "Body of a text"
            },
          },
        },
        "label": {
          "type": "string",
          "description": "Label for an instance. Can be empty if this is
            not validation data"
        },
        "metadata": {
          "type": "object",
          "description": "Any additional metadata. Can be empty field"
        },
      },
    },
  },
}
```

Figure 4.1: MongoDb JSON schema visualization

4.2 Machine Learning Component

Machine Learning (ML) Component is fully implemented in Python with the usage of open source libraries. In order to understand how to use the models, it is possible to find the code and its usage in Jupyter notebook that are stored in showcase folder. Showcase folder has four Jupyter notebooks that show both how to get the data for the models from the database and how to start using the models.

4.2.1 Data Transformers

Before models training and testing user has to fit the data transformer that transforms text into tensors form. Tensors are used as input values for models. As mentioned in chapter 2, we are working only with TF-IDF and Fast Text text encodings. Both TF-IDF and Fast Text models are created from the documents that are given to the transformer.

4.2.1.1 TF-IDF Transformer

TF-IDF transformer represents one article as a vector. From the articles are extracted all words that exist in the vocabulary and then for a document is calculated TF-IDF encoding. As a result, if we make TF-IDF encoding of a set of articles, we will get a matrix where each row represents specific document and each column represents word from a dictionary.

4.2.1.2 Fast Text Transformer

Fast Text model is a pretrained model that consists of one million words mapped to vectors. These words are stored in MongoDb. When a user starts to use Fast Text model, ML component creates a words' vocabulary from the texts taken for model training/testing. This vocabulary is created in a form of a hash map (word -> vector) where word embeddings are downloaded from MongoDb. It is important to remember that Fast Text encoding represents each word as a vector with predefined number of components. We are using word embeddings that represent each word with 300 float values. We introduce article encoding as a mean value through all words from a text that is given for encoding. Thus, if we make Fast Text encoding of a set of articles, we will get a matrix where each row represents specific document. Huge advantage of this method in comparison to TF-IDF, is that we are working only with 300 float values (300 columns if we provide encoding of set of articles) than with huge dictionary. Therefore we get lower features dimensionality and better context understanding.

4.2.2 Machine Learning

In this work we created a Generalized Model thats unifies all models. Generalized Model allows to work with each Machine Learning algorithm in same way. Generalized model is able to take a data transformer as an input argument. This feature makes it easier to work with models. In first chapter of this theses we introduced our models the way that we want to sample from their parameters' distributions. That means that we are aiming to work with ensembles. In figure 4.2 is shown a generalized diagram of machine learning structure.

We used scikit-learn [10] for basic algorithms such as Random Forests and SVMs. However, core of this project is constructed around Neural Networks. We used PyTorch [9] as a Neural Networks framework. In this work we have implemented and tested five classification algorithms. Three of them, such as SVM, Random Forest and Feed Forward Neural Network ensembles are trained on a randomly chosen subsets from the training data. For each ensemble are randomly chosen 80% from training data that are used for training. Two algorithms such as SGLD and DENFI are using full training dataset for
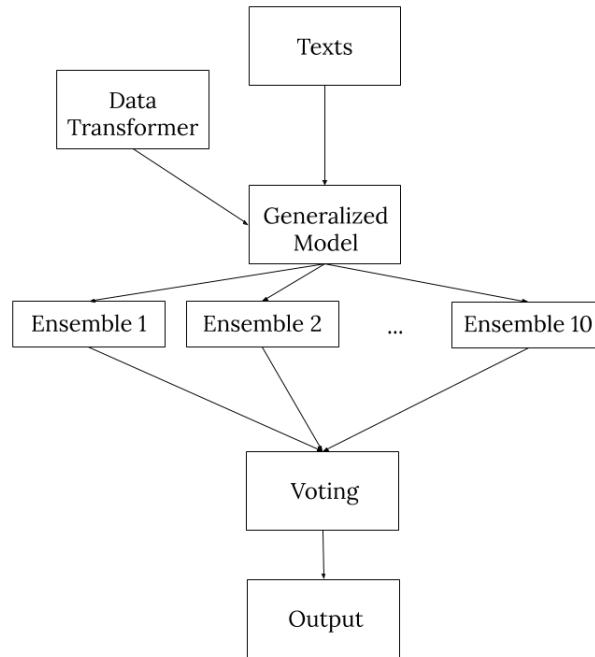
Figure 4.2: Machine Learning Workflow

their ensembles. The variability in SGLD and DENFI ensembles is reached though adding gaussian noise while ensembles training. We hardcoded amount of ensembles for all models to 10.

4.2.2.1 SVM and Random Forest Ensemble Setup

Both SVM and Random Forest models were taken and used out of the box. We created 10 SVM and 10 Random Forest ensembles with default scikit-learn setting. No modifications were provided.

4.2.2.2 Feed Forward Neural Network

Despite the fact that we used very simple Neural Networks architecture, it showed very good results. We implemented Feed Forward Neural Network with the usage of PyTorch python package with one hidden layer that consists of 100 neurons with sigmoid activation function. We chose softmax activation function for output layer.

4.2.2.3 SGLD

For SGLD we used same configuration as for Feed Forward Neural Network. One significant difference is that we used whole training dataset and were adding gaussian noise to a gradient descent. Precise configuration of SGLD can be found in github project here `https://github.com/sahanmar/Peony/blob/master/Peony_project/Peony_box/src/peony_adjusted_models/sgld_nn.py`

4.2.2.4 DENFI

Chapter 5

Passive Learning Classification

Before we can start active learning part we have to understand if implemented algorithms are capable to solve the classification task. Thus, we decided to test the models on Sports and Comedy categories from HuffPost Dataset and on Tweets from Tweets Dataset. The classification was done both for TF-IDF and Fast Text Encodings. We separated experiments with respect to the dataset types.

5.1 Passive Learning HuffPost Dataset

In this section we are illustrating ROC and AUC metrics with respect to 10−fold cross validation and 500 Sports, 500 Comedy articles. Vocabulary, that is created from 1000 articles corpus consists of 20 thousand unique words. We would like to mention that all algorithms are trained and tested with respect to 10 ensembles. Moreover, the ratio of randomly chosen training data for ensembles (SMV, Random Forest, Feed Forward NN ensembles) is set to 80%.

5.1.1 SVM Ensembles

Results for SVM Ensembles shown in figures 5.1 and 5.2. These are results both for TF-IDF and Fast Text encodings. As seen on these plots ROC and AUC values of 10−fold cross validation are extremely high. This means that our model works very good. Another interesting point is that standard deviation with respect to all runs is very low. It means that SVM ensembles could linearly separate Sports and Comedy sets with acceptable classification error.

It is also seen that ROC and AUC metrics are almost same for TF-IDF and Fast Text encodings. However, there is one significant difference. The difference is in computational time because Fast Text encoded text document consists of 300 float components and TF-IDF encoded text document consists of 20 thousand float components. Even though we are using algorithms for sparse matrix computation for TF-IDF method, computations for Fast Text encoding are five times faster. The higher amount of unique words is, the higher elapsed time per article for TF-IDF based encoding algorithm will be. Another good aspect is that because of algorithm's simplicity, SVM is trained much faster than Neural Network based models.

It is possible to conclude that the model shows good results for this classification problem and can be used for active learning experiments.
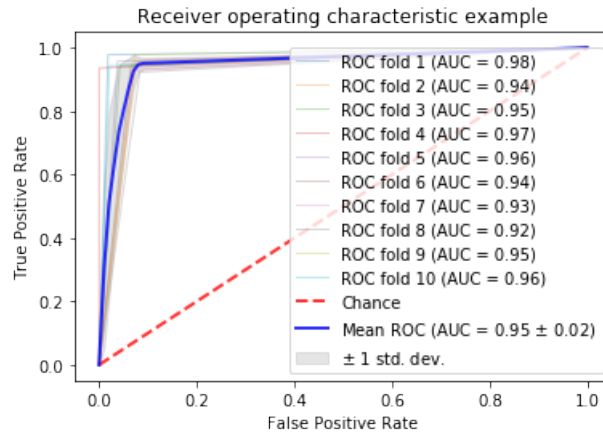
Figure 5.1: TF-IDF ROC and AUC for 10 SVM Ensembles trained and tested on Sports and Comedy data where each ensemble is trained on 80% of randomly chosen training data
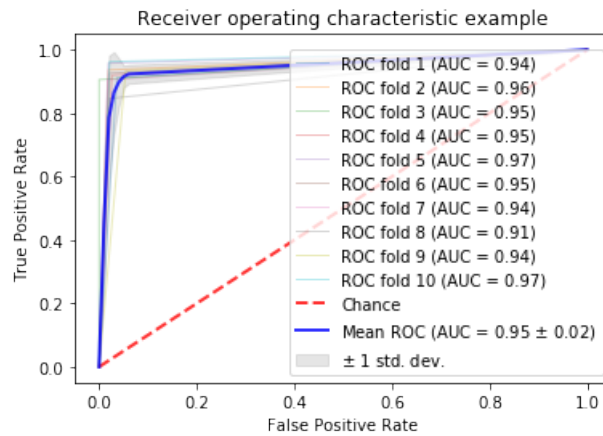

Figure 5.2: Fast Text ROC and AUC for 10 SVM Ensembles trained and tested on Sports and Comedy data where each ensemble is trained on 80% of randomly chosen training data
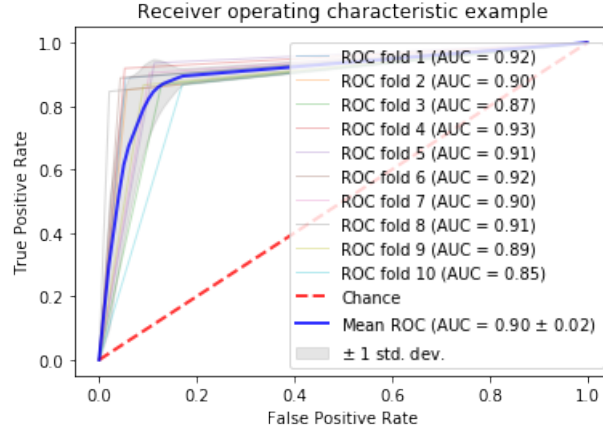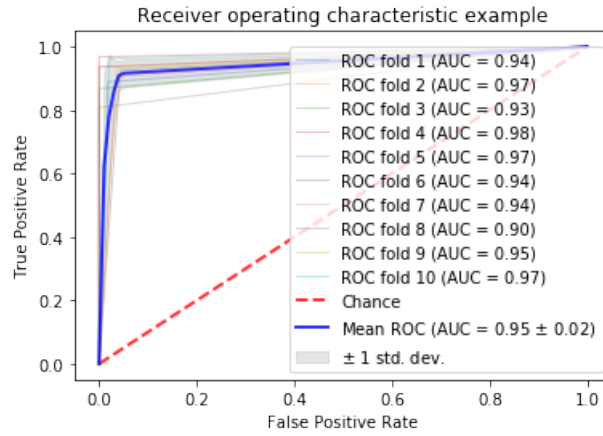
Figure 5.3: TF-IDF ROC and AUC for 10 Random Forest Ensembles trained and tested on Sports and Comedy data where each ensemble is trained on 80% of randomly chosen training data

Figure 5.4: Fast Text ROC and AUC for 10 Random Forest Ensembles Ensembles trained and tested on Sports and Comedy data where each ensemble is trained on 80% of randomly chosen training data

5.1.2 Random Forest Ensembles

Results for Random Forest Ensembles shown in figures 5.3 and 5.4. These are results both for TF-IDF and Fast Text encodings. Same as in SVM section, ROC and AUC values of 10−fold cross validation for Random Forest are high as well. This means that our model works very good. Standard deviation with respect to all runs is also very low.

If we compare results for TF-IDF and Fast Text encodings, it is seen that Fast Text based model outperforms results with respect to TF-IDF encoding model. Mean AUC value for Fast Text article encoding model is higher by 5%. It is interesting that if we apply sparse matrix algorithms for TF-IDF in Random Forest ensembles model, then computational speed for Fast Text and TF-IDF text will approximately same. This observation was made on the basis of processing 1000 articles from Comedy and Sports sets. As mentioned in SVM section we can also say that Random Forest algorithm is trained much faster than Neural Network based models.

It is possible to conclude that the model shows good results for this classification problem and can be used for active learning experiments.
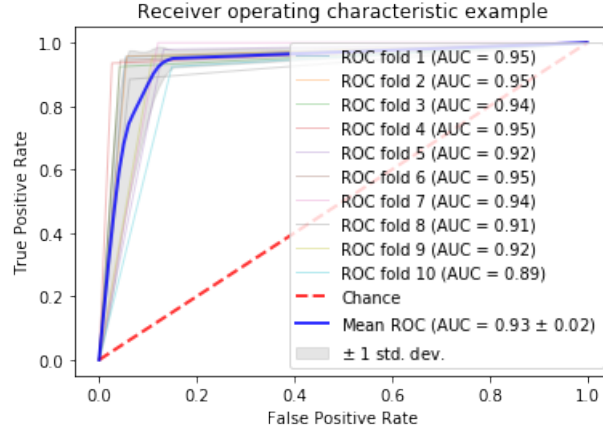
Figure 5.5: TF-IDF ROC and AUC for 10 Neural Networks Ensembles trained and tested on Sports and Comedy data where each ensemble is trained on 80% of randomly chosen training data
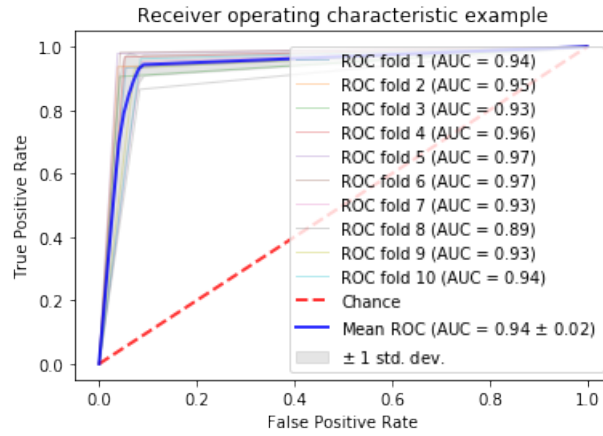

Figure 5.6: Fast Text ROC and AUC for 10 Neural Networks Ensembles trained and tested on Sports and Comedy data where each ensemble is trained on 80% of randomly chosen training data

5.1.3   Feed Forward Neural Network Ensembles

In this work we have implemented thee algorithms based on neural networks such as Neural Networks Ensemble, SGLD and DENFI algorithms. The difference between these methods is in parameters sampling. Neural Networks Ensemble takes a randomly selected subset from training data for each ensemble. Thus, we decided do not show SGLD and DENFI results in this section because they use whole set of training data. As a result if Neural Networks Ensemble shows good results, we assume that SGLD and DENFI will also perform good passive learning results.

Results for Neural Networks Ensembles shown in figures 5.5 and 5.6. These are results both for TF-IDF and Fast Text encodings. As seen in previous section, ROC and AUC values of 10−fold cross validation for Neural Networks are also high.

Comparing the results between TF-IDF and Fast Text encoding based models we can observe that Fast Text based model gives slightly better results. We have already discussed fastness of algorithm training with respect to different embedding models in Random Forest section. In case of Neural Networks, this difference is even more significant. Model that is based on Fast Text word embeddings takes

20 times less time than TF-IDF encoding based model.

5.2 Conclusion

All of tested models gave really good results in classification of Sports and Comedy categories. It means that we are able to continue working with these algorithms in active learning section. We would like to highlight that Fast Text encoding based algorithms showed a bit better results than TF-IDF. Fast Text based algorithms also showed significant improvement in fastness of algorithm trainings.

Chapter 6

Active Learning Classification

In Passive Learning section we showed that implemented algorithms are good for solving text classification task. Active learning section results is the main part of this theses. Therefor, we tested all five algorithms on the data mentioned in the Data section. The results are shown and described in further subsections. However, before starting with text classification results, we would like to show how our models represent uncertainty. As written in theoretical introduction to active learning, we use models' uncertainty for active learning loop.

6.1 Active Learning Models' Uncertainty

Due to the fact that text encoding features space dimensionality is extremely high, we introduce a 2 dimensional toy problem for uncertainty visualization. In figure 6.1 is shown a dataset that is used for toy problem classification. We generated this dataset with adding gaussian noise in order to make the task similar to real world problems.

We generated 1000 data points where 500 are assigned to Class 0 and another 500 points are assigned to Class 1. We used 50% randomly chosen data samples as a training dataset. Next step, was creating a two dimensional grid that will be used for model predictions. We assign data sample to Class 1 if prediction value is higher than 0.5. If prediction value is lower than 0.5 than the value is assigned to Class 0. We consider that model is uncertain about specific data sample if its prediction values is close
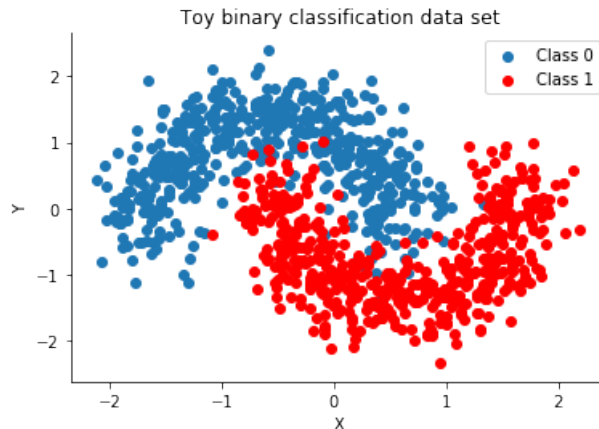

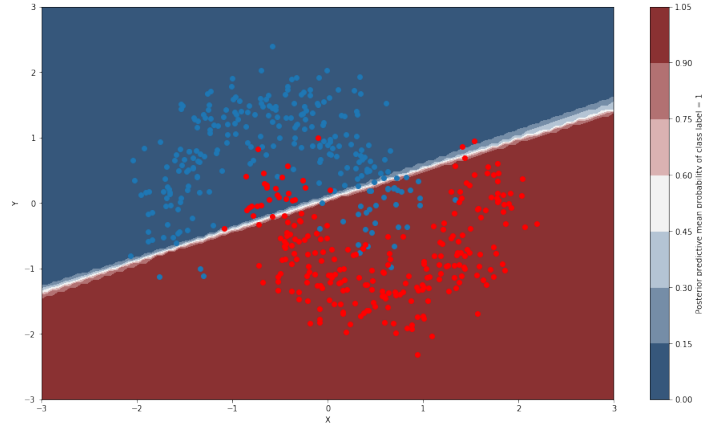Figure 6.1: Toy problem dataset visualization

Figure 6.2: SVM Ensembles posterior predictive mean probability of Class 1

to 0.5. As a result, sampling values from maximal uncertainty region will bring maximum information about the dataset. Models set up for Toy Problem is same as it is defined in Machine Learning Component section.

6.1.1 SVM Ensembles

In figure 6.2 is visualized uncertainty for SVM Ensembles model. Uncertainty bounds are linear and quite narrow. Linearity of uncertainty bounds is explained with the fact that we are using SVMs with linear kernel. Narrowness can be explained with richness of the training dataset and linear limitations of SVM decision bound.

6.1.2 Random Forest Ensembles

In figure 6.3 is visualized uncertainty for Random Forest Ensembles model. In case of Random Forest Ensembles we can see that uncertainty bounds are not linear and lay near the region where two classes are splitted. We see that uncertainty region is becoming wider near the places where datapoint of two different classes lay near each other. This is a behavior that we wanted to observe.

It is also seen that the curve is not smooth enough. This is caused with Random Forest algorithm specification.

6.1.3 Neural Networks Ensembles

In figure 6.4 is visualized uncertainty for Neural Network Ensembles model. In comparison to non-neural networks based methods that were shown above, uncertainty bounds are quite smooth. We can also observe that uncertainty bounds become wider when they go to further from the data points. This behavior can be explained with the fact that in these places algorithm did not get any training data samples. We could also see this behavior in SWM Ensembles but in this case, Neural Network models represent the uncertainties much better.

6.1.4 SGLD

In figure 6.5 is visualized uncertainty for SGLD model. We see that SGLD algorithm has similar behavior to Neural Network Ensembles. The difference is that all uncertainty bound curves has similar
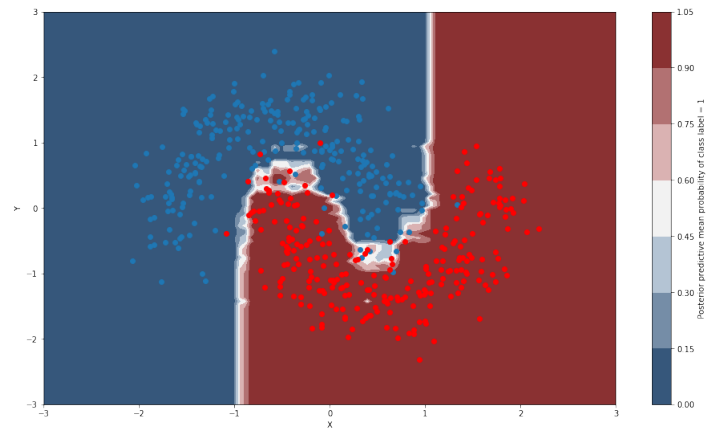
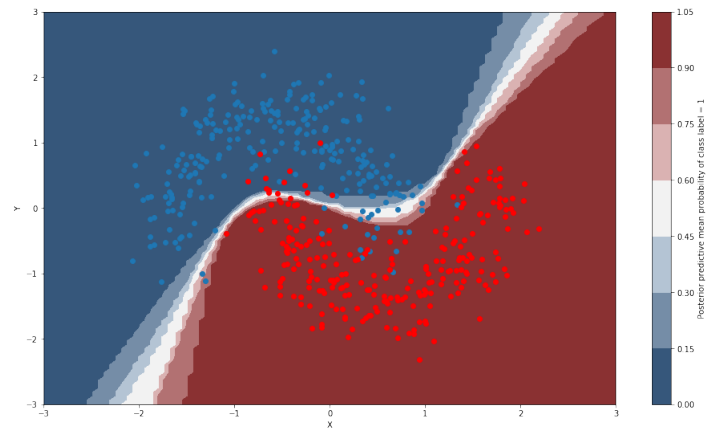Figure 6.3: Random Forest Ensembles posterior predictive mean probability of Class 1

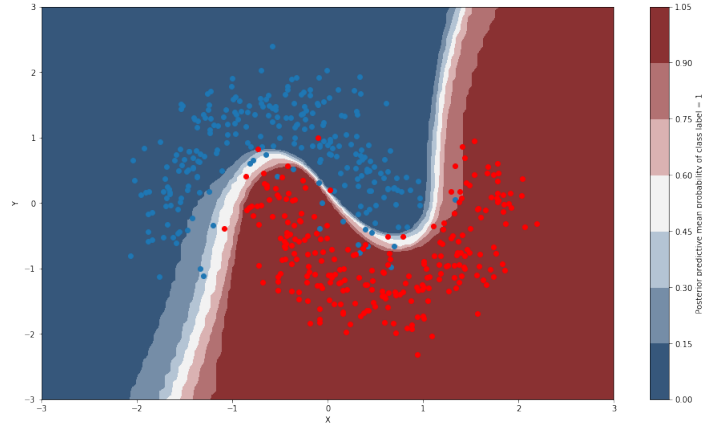Figure 6.4: Neural Network Ensembles posterior predictive mean probability of Class 1

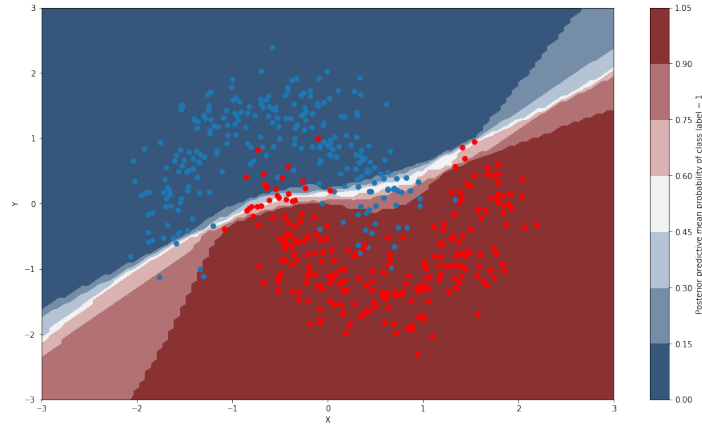Figure 6.5: SGLD posterior predictive mean probability of Class 1

Figure 6.6: DENFI posterior predictive mean probability of Class 1

curvature. This is happening due to the fact that SGLD finds loss function minimum and than samples parameters' values in a neighborhood of the minimum. As a result we expect decision bound for each parameters sample be similar to each other.

6.1.5   DENFI

In figure 6.6 is visualized uncertainty for DENFI model. We see that uncertainty bounds are very similar to Neural Network Ensembles algorithms but still a bit different. As told in Pseudocode , algorithm founds different local minimums and then we add some gaussian noise to parameters values. In next learning iterations the algorithm continues training using the weights from previous step.

The last 100 loss function values with respect to each DENFI ensemble are shown in figure 6.7. It is seen that each ensemble found its own local minimum. Additional gaussian noise adds more variability to the algorithm that makes samples more diverse.

6.1.6   Conclusion

To conclude, we can say that the best variability representation was seen in Neural Network models. Of-course, we could test SVM algorithm with the usage of different kernels but our prior interest was
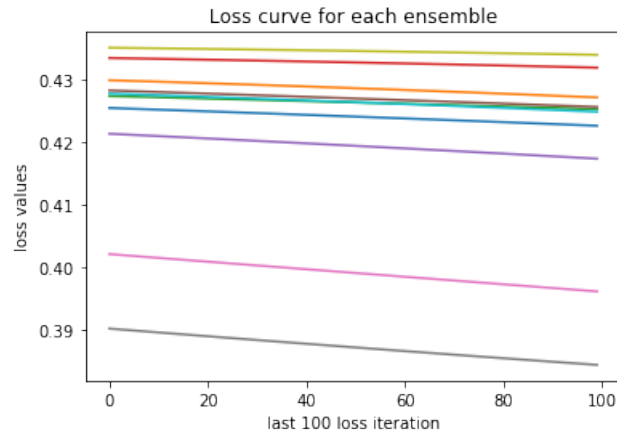
Figure 6.7: The last 100 loss function values with respect to each DENFI ensemble

around Neural Networks. We have showed that all models are able to represent variability and can be used for further tests.

6.2 Active Learning on Texts with TF-IDF Encoding Based Models

6.2.1 SVM Ensembles

6.2.2 Random Forest Ensembles

6.2.3 Neural Networks Ensembles

6.2.4 SGLD

6.2.5 DENFI

6.3 Active Learning on Texts with Fast Text Encoding Based Models

6.3.1 SVM Ensembles

6.3.2 Random Forest Ensembles

6.3.3 Neural Networks Ensembles

6.3.4 SGLD

6.3.5 DENFI

Conclusion

Text of the conclusion. . .

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