GROWING NEURAL NETWORK WITH SHARED PARAMETER

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ABSTRACT

We propose a general method for growing neural network with shared parameter by matching trained network to new input. By leveraging Hoeffding’s inequality, we provide a theoretical base for improving performance by adding subnetwork to existing network. With the theoretical base of adding new subnetwork, we implement a matching method to apply trained subnetwork of existing network to new input. Our method has shown the ability to improve performance with higher parameter efficiency. It can also be applied to trans-task case and realize transfer learning by changing the combination of subnetworks without training on new task.

1 INTRODUCTION

Deep neural networks have made great progress and shown impressive empirical results recently. But most neural networks are trained in a fixed structure which limits the expandability of neural networks. To make neural networks expandable, there is a set of methods to optimize neural network structure dynamically. It has been shown that accurate and efficient neural network can be gradually grown from a relatively small network. (Liu et al. (2019); Wang et al. (2019); Wu et al. (2021))

The existing network growing methods grow networks by the indication of loss and grow neural network by adding new parameters. In this work, we propose a new method to grow neural networks with shared parameters based on statistic theory. The key problem of growing networks is that: under what condition, the added parameters is contributive to final results. We solve this problem by combing Hoeffding’s inequality and loss descendent in the case of summation ensemble, and this qualification method is expandable to multiple types of neural network. In the implementation, we apply the proposed qualification method to neural additive model (NAM) (Agarwal et al. (2020))
which is the additive combination of multiple independent multi-layer perceptron (MLP), and grow neural network by matching trained MLP to new inputs. The general process is shown in Figure 1. Our growing method can also be applied in trans-task condition by matching trained MLP to the inputs of a new task as shown in Figure 2.

![Figure 2: Process of trans-task growing. We grow network by applying trained MLP to input of other tasks.](image)

Our growing method is a general method and can be applied in additive ensemble networks on classification tasks. Theoretical base of MLP qualification guarantees a steady improvement on performance and matching trained MLP makes higher parameter efficiency. Our growing method can also be used under trans-task condition without training and without the requirement of shared knowledge between tasks. Experiment results have shown the steady increment of performance, higher parameter efficiency in growing existing network. We also show the ability for trans-task learning, which means that function of one neural network can be altered by perturbing parameters.

2 METHODOLOGY

2.1 SUBNETWORK QUALIFICATION

For a typical neural network on a classification task, the classification layer which is normally a fully-connected layer as shown in (1),

$$ y_i = \sum_{j=1}^{m} w_{ij} f(x)_j + b_i, $$

where $f(x)_j$ is the output of former layers and $f$ could be any form of neural networks, $w_{ij}$ is the $ij$-th element of weight, $b_i$ is the $i$-th element of bias. In the rest part of this work, we refer $w_{ij} f(x)_j$ as the class-output of the $i$-th class, $y_i$ is the sum of the $i$-th class-output.

From (1), we consider that class-output for each class is a linear combination of different values. Under this condition, we grow neural network by adding new subnetwork which add new value $w_{ij} f(x)_j$ to this linear combination. In this work, to make this problem easier, each new subnetwork only focuses on one class and the class-output of other classes is set to 0.

The class which is focused is referred as target class and samples of target class is referred as target samples, other classes is referred as non-target class and samples of non-target class is referred as non-target samples.

We evaluate the grown network on a dataset which is referred as selection set. The $j$-th sample of class- $i$ is represented as $sp_{ij}$, class-output of $k$-th subnetwork on sample $sp_{ij}$ is represented as $y_{ij}^k$. 
which is equal to \( w_{ij} f(x) \) in (1). Number of classes is represented as \( N_c \), number of subnetworks is represented as \( N_n \).

The evaluation considers two aspects, from loss and from Hoeffding’s inequality. Conditions for qualifying added subnetworks is summarized at the end of this section.

2.1.1 Hoeffding’s Inequality Aspect

Hoeffding’s inequality is depicted as follows:

**Lemma 1.** Let \( X_1, X_2, \ldots, X_n \) be independent random variables such that \( a_i \leq X_i \leq b_i \) almost surely. The sum of these random variables \( S_n = \sum_i X_i \) satisfies the condition that: for all \( t > 0 \),

\[
P(|S_n - E(S_n)| \geq t) \leq 2\exp\left(-\frac{2t^2}{\sum (b_i - a_i)^2}\right)
\]

The variable \( E(S_n) \) is the expectation of \( S_n \). The estimated sum converges to the expectation of sum by probability.

From (1), since the output of each class is the sum of class-outputs of all subnetworks, Hoeffding’s inequality can be applied to our proposed growing method.

In the phase of prediction, class output of each subnetwork can be considered as a random variable. For subnetworks focus on target class \( c_t \), we use \( y_{ct,j} \) to represent the sum of class-output values on \( j \)-th sample of target-class \( c_t \) and \( y_{cnt,j} \) to represent the sum of class-output on \( j \)-th sample of non-target class \( c_{nt} \). The expectation of \( y_{ct,j} \) is represented as \( E(y_{ct,j}) \) and the expectation of \( y_{cnt,j} \) of other class is represented as \( E(y_{cnt,j}) \). If \( E(y_{ct,j}) > E(y_{cnt,j}) \), according to Hoeffding’s inequality, with enough subnetwork outputs \( y_{ct,j} \) bounded by \( [l_k, u_k] \), \( y_{ct,j} \) will converge to \( E(y_{ct,j}) \) and \( y_{cnt,j} \) will converge to \( E(y_{cnt,j}) \) by probability as shown in (2).

\[
\begin{align*}
P(|y_{ct,j} - E(y_{ct,j})| \geq t) &\leq 2\exp\left(-\frac{2t^2}{\sum (u_k - l_k)^2}\right) \\
P(|y_{cnt,j} - E(y_{cnt,j})| \geq t) &\leq 2\exp\left(-\frac{2t^2}{\sum (u_k - l_k)^2}\right).
\end{align*}
\]

Thus, the probability that \( y_{cnt,j} > y_{ct,j} \) will converge to 0, there will be a clear boundary between target samples and non-target samples.

The prediction phase is operated on a test set. In the selection phase, we evaluate added subnetworks on the selection set. We firstly suppose that the selection set and test set are formed by the same distribution. Therefore, the expectations \( E(y_{ct,j}) \) and \( E(y_{cnt,j}) \) on selection set and test set should be identical. Secondly, from Hoeffding’s inequality, with enough subnetworks, there will be a clear boundary between target samples and non-target samples. From the above supposition, we can select subnetworks with the following conditions.

\[
\sum_{ij} w_{ij}^k y_{ij}^k > 0, \quad w_{ij}^k = \begin{cases} 
\max(\max(y_{cnt,j}) - y_{ij}, 0), & i = c_t \\
\min(\min(y_{ij}) - y_{ij}, 0), & i \neq c_t.
\end{cases}
\]

In (3), the average class-output value on target samples should be greater than the average class-output values on non-target samples, this condition guarantees \( y_{ct,j} > y_{cnt,j} \). In (4), we use a weighted sum of class-output on each sample, to make \( \min(y_{ct,j}) > \max(y_{cnt,j}) \).

2.1.2 Loss Aspect

From loss aspect, purpose of added subnetworks is to minimize loss. A new subnetwork will add value on both target samples and non-target samples. We decide which sample should be added more value and which sample should be added less value.

For cross-entropy loss, we compute derivative of loss by \( y_{ij}^k \), as shown in (5) and (6).

\[
- \frac{\partial l_j}{\partial y_{ct,j}^k} = 1 - \frac{1}{\tau_j + 1}, \quad \tau_j = \frac{\sum_{i \neq c_t} e^{y_{ij}^k}}{e^{y_{ct,j}^k} + \sum_{i \neq c_t} e^{y_{ij}^k}},
\]
\[
\frac{\partial l_j}{\partial y_{c,ij}^k} = \begin{cases} 
1 & \tau_j = \sum_{i \neq c_r} \frac{e^{y_{ij}}}{e^{y_{ij}} + e^{y_{c,ij}^k}}, \ c_r \neq c_t, \\
0 & \end{cases}
\]

In (5), if the true label is target class, for the same positive \(y_{c,ij}^k\), the greater \(\tau_j\) is, the greater loss descendent. In (6), if the true label is non-target class, the greater \(\tau_j\) is, the smaller loss increment. Therefore, purpose of added subnetwork is to make \(\tau_j\) as great as possible on each sample.

The ideal condition of minimizing loss is that output values on all target samples is identical as \(\frac{y_{c,ij}^k}{y_{c,nt,j}^k}\) and all output values on all non-target samples is identical as \(\frac{y_{c,nt,j}^k}{y_{c,nt,j}^k}\), and the variance should be minimized to 0. This condition also satisfies the condition in (4). Thus, the derivative of variance on target class and non-target class is shown in (7) and (8) respectively.

\[
- \frac{\partial \text{Var}(y_{c,ij})}{y_{c,ij}} = \frac{2}{N_c} (E(y_{c,ij} - y_{c,ij})),
\]

\[
- \frac{\partial \text{Var}(y_{c,nt,j})}{y_{c,nt,j}} = \frac{2}{N_{c,nt}} (E(y_{c,nt,j} - y_{c,nt,j})).
\]

In conclusion, in a fixed dataset, for one added subnetwork which focus on one target class \(c_t\). This added subnetwork should make more contribution to target samples and make contribution to lessen variance as depicted in (9).

\[
\begin{align*}
\text{If } & y_{c,ij}^k > y_{c,nt,j}^k & \sum_{ij} w_{ij}^k y_{ij}^k > 0, \\
\text{we } & & w_{ij}^k = \begin{cases} 
E(y_{c,ij} - y_{c,ij}), & i = c_t \\
E(y_{c,nt,j} - y_{c,nt,j}), & i \neq c_t
\end{cases}
\end{align*}
\]

2.1.3 ELECTION

If we use a threshold \(\text{thd}_k\) to each added subnetwork output value as (10),

\[
y_{ij}^k = \begin{cases} 
1, & y_{ij}^k > \text{thd}_k \\
0, & y_{ij}^k \leq \text{thd}_k
\end{cases}
\]

we can leverage the form of Hoeffding’s inequality on binary distribution as shown in (11) and (12).

\[
P(y_{c,ij} \leq (p_{ij} - \varepsilon)N_n) \leq e^{-2\varepsilon^2 N_n} \quad 0 < \varepsilon < p_{ij},
\]

\[
P(y_{c,nt,j} \geq (p_{nt,j} + \varepsilon)N_n) \leq e^{-2\varepsilon^2 N_n} \quad 0 < \varepsilon < 1 - p_{nt,j}.
\]

\(y_{c,ij}\) is the sum of class-output on target sample \(s_{y_{c,ij}}\), \(y_{c,nt,j}\) is the sum of class-output on non-target sample \(s_{y_{c,nt,j}}\). \(p_{ij}\) is the probability that class-output on \(s_{y_{c,ij}}\) is 1 and \(p_{nt,j}\) is probability that class-output on \(s_{y_{c,nt,j}}\) is 1. If \(\min(p_{ij}) > \max(p_{nt,j})\) and with enough subnetwork, \(y_{c,ij}\) will converge to \(p_{ij}N_n\) and \(y_{c,nt,j}\) will converge to \(p_{nt,j}N_n\) by probability. Therefore, will form a clear boundary between target samples and non-target samples.

We suppose that samples of each class are equal in selection set, we alter (9) to (13) as follows.

\[
\begin{align*}
\text{If } & \text{prec}_k > 1/N_c & \sum_{ij} w_{ij}^k y_{ij}^k > 0, \\
\text{we } & & w_{ij}^k = \begin{cases} 
E(y_{c,ij} - y_{c,ij}), & i = c_t \\
E(y_{c,nt,j} - y_{c,nt,j}), & i \neq c_t
\end{cases}
\end{align*}
\]

\(\text{prec}_k\) is the precision of target samples in all the samples on which class-output is set to 1.

We use this principle for election by subnetworks in trans-task case and experiment show that election mode is able to improve performance steadily without training new parameters. Detailed deduction on election is shown in Appendix A.

2.2 BASE NETWORK

We choose NAM as the base network to verify our growing method. NAM consist of multiple independent MLPs which focus on specific range of input. Output of MLP is summed as the class-output for classification. NAM is formulated as (14),

\[
y_i = \sum_k^n f(x_k)i,
\]

(14)
where $f(x)$ is the function of MLP, $x_k$ is the input of the $k$-th MLP. Since NAM is tree structure network, we refer each MLP as branch in the rest part of this paper. In the setting of NAM, branch represent the subnetwork in section 2.1.

In this work, we firstly train an NAM as base network. Our growing method is to apply trained branches to other inputs to add branches to base network. In trans-task condition, we apply trained branches of one task to inputs of a new task to grow network in the new task.

2.3 GROW WITH SHARED PARAMETERS

Since we know the condition for qualified branches, we can try different branch to one input range and select qualified branches on a selection set. But randomly trying branches is not efficient in computation. In this part, we solve this problem by matching distributions.

Based on theories of branch qualification, the ideal condition for matching a trained branch to one input range is that: we select one target-class $c_t$ on one input range and apply a trained branch to this input range, and the class-output $y^b_{ij}$ satisfies the condition of (9). Namely, class-output on samples of class $c_t$ is higher than class-output of samples of other class.

The problem of matching input ranges and branches is solved by matching distributions. For original trained branches, the structure and input-output distribution of each class are shown in Figure 3, we randomly select two dimensions from the input space, and each point in the input space is an input sample. In the distribution heat map, input sample is randomly generated from uniform distribution. For each class, some part of input space obtains higher output value and other part obtains lower output value, namely, each class focuses on specific part of input space. For input range to be matched, we select one class as target class and other class as non-target class, the distribution is shown in the left scatter map of Figure 4. Our purpose for matching is to find a candidate branch, target samples matches the branch’s high-value part as shown in Figure 4, and make class-output on target samples higher than class-output on other samples. The matched branch only focuses on the target class, namely, the red point in Figure 4.

To make a clear definition, input range to be matched is referred as reference input, sample of reference input is referred as reference sample. The class to be focused on reference input is referred as target class $c_t$, other class is referred as non-target class. Branch to be matched is referred as candidate branch, class which is selected for class-output is referred as branch-class $c_b$.

With the theoretical base of branch qualification and the purpose of matching branches, the overall process is summarized in Figure 5. We grow neural network in an iterative fashion. In the clustering step, we transfer the distribution of high value part in Figure 2 to a set of clusters. In the matching step, we match target sample to clusters of high value part and make target samples obtain higher class-output values. In the branch selection step, we select qualified branches by (9) and (13). in the steps of adding network, we show how to add branches to current network for further tuning or election.
2.3.1 Clustering

We consider the high value part in Figure 3 as the combination of multiple clusters, and clustering method convert high-value distribution to combination of clusters.

We do cluster on each class in an iterative way. First, we use multi-variate Gaussian kernel to find the densest point in the branch input space. Then we use nearest neighbor method to search all the samples of this cluster and form a cluster. Finally, we remove the selected samples and repeat the former steps to find new clusters.

Each cluster is a pair contains class $c_{bi}$, cluster center $sp_{bi}$ max class-output value $y_{cbi max}$ as $\{c_{bi}, sp_{bi}, y_{cbi max}\}$. We depict a cluster as a center with its maximum class-output. Detailed implementation of clustering is shown in Appendix B.

2.3.2 Matching

With branch clusters, we can match reference samples to branch clusters. We assume if one reference input point is closer to the center of one cluster, the class-output will be greater.

Based on this supposition, we can evaluate matching by distance. For branch class $c_{bi}$ with multiple cluster $Cl_{cbi}$ and reference class $c_{ri}$ with $N$ reference samples $sp_{ri}$, partial average distance is defined in (15),

$$d_{b_{ri}} = \frac{1}{n} \sum_{j} w_{nst} \sqrt{(sp_{c_{ri}j} - Cl_{nst})^2}, \quad w_{nst} = \frac{e^{y_{nst max}}}{\sum_{j} e^{y_{nst max}}}, \quad n \leq N. \quad (15)$$

$Cl_{nst}$ is the nearest center from $sp_{c_{ri}j}, y_{nst max}$ is the max class-output value of the nearest cluster, $n \leq N$ means if there are $n$ sample satisfies the matching condition that $d_{b_{ri}}$ is the smallest distance among all reference classes, we set a boundary in input space and drop samples out of boundary.
We iteratively evaluate combinations of all $c_b$ and $c_r$, and find the best matched reference class for each branch class.

After finding the best-matched candidate branch and corresponding branch class $c_b$ on one reference class $c_r$, we add this branch to candidate branches for branch selection step, and the reference class $c_r$ is marked as target-class $c_t$ of new branch. Only class output of $c_b$-th class will be used as class output of the matched branch and class-output of other branch class is set to 0.

We also developed parameter transfer method to make matching easier as shown in Appendix C.

2.3.3 BRANCH SELECTION

Matched branches will be evaluated in a selection dataset, which is a subset randomly sampled from training dataset. The qualified branches satisfy the conditions in (9) and (13) in branch qualification part.

For the qualification condition, we use (13) to check each added branch. We also use a threshold value to select class-output values greater than this threshold. Class-output values lower than threshold are set to 0. To make the contributions of target class greater than non-target class, there will be a non-negative bias added to all the above-threshold samples in the tuning step.

Qualified added branches will be used on further tuning or election.

2.3.4 TUNING

After branch qualification, we add a class-mask which is formulated in (16) to each added branch.

$$
y_{ij}^k = \text{ReLU}(a) \left( \text{ReLU} \left( \frac{y_{ij}^k - \text{thd}}{v_{span}} \right) + s(y_{ij}^k) \text{ReLU}(b) \right), \quad s(y_{ij}^k) = \begin{cases} 
1, & y_{ij}^k > \text{thd} \\
0, & y_{ij}^k \leq \text{thd}
\end{cases}
$$

(16)

the variable $v_{span}$ is the value span $\max(y_{ij}^k) - \text{thd}$ computed in the selection set and applied in tuning and testing, $a$ and $b$ are trainable variable as scale and bias, scale and bias are set to positive using ReLU. This setting is to make each branch make more contribution to target samples than non-target samples. As shown in (16), only scale and bias will be trained in the tuning mode, and once class mask is trained, parameters is fixed and won’t be trained later, this makes tuning efficient.

2.3.5 ELECTION

In the election mode, we compute class score as (17).

$$
\text{score}_{ij} = \sum_k y_{ij}^k - E(y_{ij}^k) / \sigma(y_{ij}^k),
$$

(17)

$E(y_{ij}^k)$ and $\sigma(y_{ij}^k)$ is the mean and std of $y_{ij}^k$. These two values are computed on training set and applied in test set. The prediction of class is the class which has the maximum score among all class scores.

3 EXPERIMENT

We test our growing method on two tasks: 1) adding new branches to base network to improve the performance of original network, we show our growing method can gradually improve performance with higher parameter efficiency. 2) We also use branches in base network to do election on another task, and we show our growing method can gradually improve performance in trans-task case. Our growing method can transfer network without training.

3.1 BASE NETWORK

We train base network on cifar10 dataset. Cifar10 dataset which contains 50000 training samples and 10000 testing samples is used for image classification task. We normalize all samples to the range of [-0.5, 0.5] and no data augmentation is applied. Each image in cifar10 is of [3, 32, 32].
Table 1: performance and parameters on cifar10

| Network                | Accuracy | Loss   | parameters |
|------------------------|----------|--------|------------|
| Base network           | 0.468    | 1.4684 | 33750      |
| Full perception network| 0.5022   | 1.3873 | 135000     |

We set each branch of base network percept the input range of \([6i : 6i + 3; 6j : 6j + 3], i, j \in \mathbb{N}\) of each channel, each branch contains 4 layers and each layer is a fully connected layer with dimension 9. Base network originally contains 75 branches. We train the base network by gradient descend method and use Adam optimizer for optimization.

3.2 Network Growth

After training base network, we apply our growing method to the base network. We add branches to all the input ranges \([i : i + 3; j : j + 3], i, j \in \mathbb{N}\). Figure 6 depicts the accuracy and loss on test set with the increment of added branches. We can see a steady increment on accuracy from 0.468 to 0.5016 and a steady decrement on loss from 1.4684 to 1.4379. To compare the parameter efficiency, we also trained a full-perception NAM which percept all the input ranges and all channels of input image as the comparison network of NAM on cifar10 task, performance and parameters of full perception network is shown in Table 1. The grown network added totally 1456 branches, each added branch contains 2 new parameters of class-mask. Therefore, training parameters of grown network is 36662, which is much less than parameters of full perception network, and the performance of grown network is comparable with full perception network.

From Figure 6, we can see a steady improvement on performance, which verified the effectiveness of our growing method. By comparing the performance of grown NAM with the full-perception NAM, the performance of grown network has shown a higher parameter efficiency.

![Figure 6: Loss and accuracy of grown network as new branches added.](image)

3.3 Transfer Task Learning

After training the base network on cifar10 dataset, we apply branches of base network to MNIST dataset. MNIST is a dataset which contains 50000 training samples and 10000 test samples and used for written number recognition. We normalize each sample to range \([-0.5, -0.05]\) and there is no data augmentation.
There is initially no branch on MNIST task, and we gradually apply the originally trained 75 branches of base network to this task, and use election mode for prediction. Figure 7 depicts the accuracy on test set with the increment of branches. From Figure 7, there is a steady increment on the performance and achieve 0.8991 on train set and 0.8813 on test set. This performance is realized by applying the original 75 branches to different input ranges of MNIST task.

The experiment indicates that: by applying branches in different input range, branches trained on one task can be applied to another task without training new parameters, and changing function of one neural network can be realized by perturbating the order of subnetworks.

To compare the representability of NAM on MNIST task, we train a full perception NAM, the performance is shown in Table 2. There is a gap of 0.08 on accuracy, which means our current trans-task growing is not better than training directly.

| Network                        | Accuracy | Loss  | parameters |
|--------------------------------|----------|-------|------------|
| Full perception network        | 0.9613   | 0.1362| 36450      |

There have been several works focus on growing neural networks to gradually improve performance on a given task. (Schwenk & Bengio (2000); Bengio et al. (2006); Bach (2017)) progressively adding new parameters to existing neural network by the indication of gradient while keeping previous parameters fixed. However, this set of works need to solve complex non-convex optimization problems. Our growing method also add new branches with previous parameters fixed, but leveraging statistic base to avoid optimization problems.

Another set of growing methods grow networks by splitting neurons, (Wynne-Jones (1991); Chen et al. (2015)) grow neural network by heuristic or random strategies, (Liu et al. (2019)) propose the Splitting Steepest Descendent which provide a new method to selectively split neurons while guarantee the improvement of performance, but the growing operation is limited to splitting neurons. (Wu et al. (2021)) further extent the growing operations to both deepen and widen neural networks and optimize the grown network in a more efficient way. These splitting methods is efficient on
growing neural network by the indication of loss. However, these methods need to introduce new training parameters to existing neural network.

Neural architecture search (NAS) methods also including adding new parameters to existing neural network. [Hu et al. (2019)] proposed an efficient way of growing neural networks by adding shortcuts as gradient boosting. In our work, we use both gradient indication and statistic indication.

5 Conclusion

In this work, we propose a method for growing network with shared parameter and the theoretical base of our growing method. Our growing method provide a possible way to dynamically improve the performance of a given network with few added parameters, and a new approach for transfer learning without tuning parameter but change the combination of different parameters. For future work, we will make growth on add new layers, do matching based on attention models, make growth more efficient and apply our growing method to more tasks. Source code of our growing method can be found in https://github.com/Rain-axt/Growing-neural-networks.

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A Appendix

detailed deduction of election:

Let $X_1, X_2, \ldots, X_n$ be independent random variables such that $a_i \leq X_i \leq b_i$ almost surely. The average value of these random variables $\overline{X} = \frac{1}{n} \sum_i X_i$ satisfies the condition that: for all $t > 0$:

$$P(\overline{X} - E(\overline{X}) \geq t) \leq \exp\left(-\frac{2n^2t^2}{\sum_i (b_i - a_i)^2}\right).$$
if each random variable is limited within 0, 1 and satisfy distribution in (A.1),

\[ P(X_i = 1) = p_1, \quad \text{(A.1)} \]

The average form of Hoeffding’s inequality can be transferred to (A.2)

\[ P(\bar{X} - n_{nt} \geq t) \leq \exp(-2nt^2). \quad \text{(A.2)} \]

If each added subnetwork output value is limited within 0, 1, given \( y_{ctj} \) is the sum of class-output of target sample \( sp_{ctj} \), \( y_{ntj} \) is the sum of class-output of non-target sample \( sp_{ntj} \), \( p_{tj} \) is the probability that class-output on \( sp_{ctj} \) is 1 and \( p_{ntj} \) is probability that class-output on \( sp_{ntj} \) is set to 1. According to (A.2), we have (11) and (12).

\[ P(y_{ctj} \leq (p_{tj} - \varepsilon)N_n) \leq e^{-2\varepsilon^2 N_n} \quad 0 < \varepsilon < p_{tj} \]
\[ P(y_{ntj} \geq (p_{ntj} + \varepsilon)N_n) \leq e^{-2\varepsilon^2 N_n} \quad 0 < \varepsilon < 1 - p_{ntj} \quad \text{(11, 12)} \]

Consider there is totally \( N_k \) samples on which class-output is set to 1 from the \( k \)-th added subnetwork, the precision of target class \( prc_t^k \) and on non-target class \( prc_{nt}^k \) is defined as (A.3),

\[ prc_t^k = \frac{N_t^k}{N_k}, \quad prc_{nt}^k = \frac{N_{nt}^k}{N_k}. \quad \text{(A.3)} \]

For all the subnetworks, the precision is defined as (A.4),

\[ prc_t = \frac{N_t}{N}, \quad prc_{nt} = \frac{N_{nt}}{N}. \quad \text{(A.4)} \]

where \( N_t \) and \( N_{nt} \) is the sum of \( N_k^t \) and \( N_k^{nt} \) respectively.

From (8) and (9), variance of class-output on target samples is minimized to 0 and class-output non target samples is also minimized to 0. In the ideal condition, precision of all non-target class will be identical. Therefore, we have the relation between \( prc_t \) and \( prc_{nt} \) as (A.5),

\[ prc_{nt} = \frac{1 - prc_t}{N_0 - 1}. \quad \text{(A.5)} \]

To make a clear boundary between target samples and non-target samples, precisions should satisfy \( prc_t > prc_{nt} \), substitute (A.5) to this condition, we have condition of \( prc_t \) as show in (A.6),

\[ prc_t > \frac{1}{N_c}. \quad \text{(A.6)} \]

We expand this condition to each subnetwork precision and keep minimizing variance of class-output on target samples and non-target samples. The final condition is depicted as (13),

\[ \begin{align*}
    prc_t^k &> 1/N_c, \\
    \sum_{ij} w_{ij}^k y_{ij}^k &> 0,
\end{align*} \]
\[ w_{ij}^k = \begin{cases} E(y_{ctij}) - y_{ctij}, & i = c_t \\
                    E(y_{ntij}) - y_{ntij}, & i \neq c_t \end{cases}. \quad \text{(13)} \]

### B Appendix

detailed implementation of clustering:

We firstly generate samples which are randomly drawn from uniform distribution for candidate branches, we compute class-output on each class and save the pairs \{sample, class, class-output\} as \( \{sp_{ij}, c_{bi}, y_{bi}j\} \), we refer these pairs as branch-pair. For each branch-class \( c_{bi} \), only pairs with top 20% class-output value \( y_{bi}j \) are reserved and used for clustering. Reference samples is selected from train-set, the labels of reference samples are the true labels \( c_{tr} \) from dataset, reference pairs \{sample, true-class\} are saved as \( \{sp_{P}, c_{tr}\} \).

In this work, we use mean shift cluster (MSC) as our clustering method. The pseudo code of clustering is shown in Table B.1. Mean shift cluster is an iterative method to shift one point to
Table 3: Pseudo code of clustering

**Algorithm 1: Clustering**

input: set of samples of class $c_i$, $S = \{s_{p_j}, c_{bi}, y_{c_{bij}}\}$; neighbor distance $d_{nb}$; minimum shift distance $d_{min}$, empty cluster set $S_{cl}$

While $|S| > 0$ do
  Random select start point $s_{p_{cr}}$
  While shift distance $> d_{min}$ do
    Compute shifted point $s_{p_{sft}}$ by (B.2)
    \[ \text{dist}(s_{p_{cr}}, s_{p_{sft}}) \rightarrow \text{shift distance} \]
    \[ bms_{p_{sft}} \rightarrow s_{p_{cr}} \]
  End
  Select neighbor samples $\{s_{p_{nb}} - \text{dist}(s_{p_{cr}}, s_{p_{sft}}) < d_{nb}\}$ as cluster
  Add cluster to $S_{cl}$
  Remove neighbor samples from $S$
End

denser point and gradually approach to the densest point. We use multi-variate Gaussian kernel to compute density weight for each sample as depicted in (B.1).

\[
G(s_{p_1}, s_{p_2}) = \frac{1}{2\pi |\text{cov}|^{-1}} \exp\left( (s_{p_1} - s_{p_2})^T \text{cov}^{-1} (s_{p_1} - s_{p_2}) \right). \tag{B.1}
\]

The shifted point is the weighted sum of all current samples as shown in (B.2).

\[
s_{p_{sft}} = \sum_j \frac{G(s_{p_{cr}}, s_{p_j})}{Z} s_{p_j}, \quad Z = \sum_j G(s_{p_{cr}}, s_{p_j}). \tag{B.2}
\]

We stop shifting point if the shifted point is close enough to current point, and current point is the densest point. After finding the densest point, we select neighbor point within a distance and form a cluster of samples.

Each class of branch pairs are clustered separately. To make clustering suitable for different distribution, we normalize branch samples to 0-mean and 1-std before clustering and do denormalization after clustering.

For each cluster, the max class-output value and corresponding sample are marked as center of this cluster, which means higher output value is near this center. We represent one cluster $C_{bi,j}$ as $\{c_{bi}, s_{p_{cr}}, y_{c_{bij}}^{max}\}$, each cluster is a pair contains class $C_{bi,j}$ cluster center $s_{p_{cr}}$, max class-output value $y_{c_{bij}}^{max}$. We depict a cluster as a center with its maximum class-output.

**C Appendix**

detailed implementation of parameter transfer:

To make matching easier, we normalize both branch clusters and reference samples. Normalization will make patterns of branch clusters and reference samples closer.

For branch clusters $C_{bi}$ of one branch class $c_{bi}$, there is multiple samples within each cluster and we use set($x_{bi}$) to represent the sample set, and each cluster is pair of class, center, max class-output $\{c_{bi}, s_{p_{cr}}, y_{c_{bij}}^{max}\}$. We normalize clusters of $c_{bi}$ as (C.1),

\[
s'_{p_{ci}} = \begin{cases} \frac{s_{p_{ci}} - \max(x_{bi})}{\max(x_{bi}) - \min(x_{bi})} [d_1, d_2, \ldots, d_m] \\ x_{bi}[d_{k_1}] \leq x_{bi}[d_{k_2}] \end{cases} \tag{C.1}
\]

The operation of $[\ ]$ is the dimension perturbation which reorder dimensions of each cluster-center. we sort each dimension to make average values increase by dimensions. $s'_{p_{ci}}$ is the normed center.
For reference samples \( sp_{r,j} \) of one reference class \( c_{ri} \), we apply the same normalization method to normalize each sample. It worth noting that, the dimension order of cluster and reference samples are not identical.

After matching candidate branch and reference input, we applied a parameter transfer step to counter-act normalization. If one point in the reference space matches the point of branch-input space as shown in (C.2),

\[
\begin{align*}
  \frac{x_{ri}[i_r] - x_{rb}[i_r]}{\sigma(x_{ri}[i_r])} &= \frac{x_{bi}[i_b] - x_{bb}[i_b]}{\sigma(x_{bi}[i_b])} \\
  \sigma(x_{ri}[i_r]) &= \max(x_{ri})[i_r] - \min(x_{ri})[i_r] \\
  \sigma(x_{bi}[i_b]) &= \max(x_{bi})[i_b] - \min(x_{bi})[i_b]
\end{align*}
\]

and the computation of the first layer of matched branch is shown in (C.3),

\[
y_{di} = \sum_{i_b} w_{di,i_b} x_{bi}[i_b] + b_{di},
\]

the purpose of parameter transfer is to find a transferred parameter to satisfy the condition in (C.4),

\[
\sum_{i_r} w'_{di,i_r} x_{ri}[i_r] + b'_d = y_{di},
\]

namely, the output of branch-input sample should be identical to output of reference sample.

We can make a transformation from \( x_{ri}[i_r] \) to \( x_{bi}[i_b] \) as (C.5).

\[
x_{bi}[i_b] = \sigma(x_{bi}[i_b]) \frac{(x_{ri}[i_r] - \overline{x_{ri}[i_r]})}{\sigma(x_{ri}[i_r])} + \overline{x_{bi}[i_b]}
\]

By substituting (C.5) into (C.3), we have (C.6).

\[
y_{di} = \sum_{i_b} w_{di,i_b} \frac{\sigma(x_{bi}[i_b])x_{ri}[i_r]}{\sigma(x_{ri}[i_r])} - \sum_{i_b} w_{di,i_b} \frac{\sigma(x_{bi}[i_b])x_{ri}[i_r]}{\sigma(x_{ri}[i_r])} + \sum_{i_b} w_{di,i_b} \overline{x_{bi}[i_b]} + b_{di}.
\]

Form (C.6), we can obtain the transfer weight and bias as (C.7).

\[
\begin{align*}
  w'_{di,i_r} &= w_{di,i_b} \frac{\sigma(x_{bi}[i_b])}{\sigma(x_{ri}[i_r])} \\
  b'_d &= b_d - \sum_{i_b} w_{di,i_b} \frac{\sigma(x_{bi}[i_b])x_{ri}[i_r]}{\sigma(x_{ri}[i_r])} + \sum_{i_b} w_{di,i_b} \overline{x_{bi}[i_b]}
\end{align*}
\]

This step of parameter transfer is not necessary for matching but will greatly improve the efficiency of matching. In practice, we apply parameter transfer to accelerate matching.