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NVIDIA Deep Speech

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Supervisor's statement

Hereby I confirm that the presented thesis was prepared under my supervision and that it fulfils the requirements for the degree of Bachelor of Computer Science.

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Authors' statements

Hereby I declare that the presented thesis was prepared by me and none of its contents was obtained by means that are against the law.

The thesis has never before been a subject of any procedure of obtaining an academic degree.

Moreover, I declare that the present version of the thesis is identical to the attached electronic version.

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Abstract

The authors of this thesis focus on implementing scripts for training DeepSpeech2 model for Automatic Speech Recognition. We try to reproduce results obtained by Baidu Research in End-to-End Speech Recognition paper [2] using `PyTorch` framework. We also experiment with obtaining dataset for Polish language and trying DeepSpeech2 model for it. Finally, we provide fully trained models for English and Polish together with statistics about how changing hyperparameters and architecture impacts model's performance and accuracy.

Keywords

Deep Speech, ASR, Neural Networks, Machine Learning, Python, PyTorch, NVIDIA, RNN, multi-GPU, FP16

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Introduction

Transcription of spoken language is a crucial problem for many areas of modern technology industry. Being able to communicate with electronic devices not only by touching them, but also by talking to them is an important goal for IT companies. Such devices are more user-friendly so it is for sure beneficial for everybody. To achieve this goal various solutions were proposed and many of then use complex algorithms (e.g. Hidden Markov Models) [1]. However, it has been shown that the best accuracy can be achieved with Automatic Speech Recognition (ASR) models based on neural networks [2].

Our thesis concentrates on implementing state of the art ASR model DeepSpeech2 described in [2] and we realize it with the support of NVIDIA Corporation. Authors of DeepSpeech2 prepared their model only for recognizing English and Mandarin, so we plan to experiment with applying it to Polish language as well. We think, it is the biggest challenge, since accuracy of the model depends not only on its implementation, but also on the size and diversity of used dataset. Therefore, we have to find appropriate one (paying attention to licenses and copyrights) and prepare it adequately. Large size of the dataset creates another problem – we need our model to be able to train on that data in reasonable time and then work in real time. Last but not least, in order to determine the best hyperparameters we have to run many experiments, collect their results, and finally analyze them.

In order to accomplish our goals we are going to implement DeepSpeech2 model using `PyTorch` deep learning framework, which is supported with CUDA, and is considered to be comfortable to work with. To achieve high performance system we plan to use open-source libraries prepared by NVIDIA which makes it possible to train one neural network over multiple GPUs. Another optimization which we expect to speed computations up is using half precision floating-point numbers (also known as FP16) instead of single precision. When it comes to collecting datasets we assume it should be relatively easy for English as there are lots of free English utterances with transcriptions. However, for spoken corpus of Polish it may be harder - we are going to search for Polish speech collected for university programs and from audiobooks.

Structure of our thesis is the following. In Chapter 1 we introduce architecture of DeepSpeech and DeepSpeech2 models in terms of, among others, used layers, data flow and functions. After that in Chapter **??** we present applied optimizations which increased network performance. Next in Chapter 2 we describe the results of experiments on model hyperparameters. In Chapter **??** we describe how we modified and trained neural network to detect Polish language and compare obtained results with English model. Finally, in Chapter 3 we summarize our experiments, show their results and present final version of the model.

We divided our work on the model into two main parts. The first one consisted of preparing appropriate datasets (for both English and Polish) and processing them to fit the model – Łukasz Kondraciuk and Jan Tabaszewski were in charge of this part. The second one consisted of implementing the model and applying GPU optimizations to it – this was the task for Piotr Ambroszczyk and Wojciech Przybyszewski.

Chapter 1

Basic model description

DeepSpeech 2 (DS2) system is a recurrent neural network trained to ingest speech spectrograms and generate a text transcription. Following description of the model architecture is pretty basic and it is based on [2]. More details can be find there. **Czy coś takiego jak wyżej w pełni wystarcza, żebyśmy mogli cytować zdania z tej pracy bez uprzedzenia?**

1.1. Input and Output specification

Let $\{(x^{(1)}, y^{(1)}), (x^{(2)}, y^{(2)}), ...\}$ be a training set. $x^{(i)}$ is a time-series of variable length where every time-slice is a spectrogram of power normalized audio clips, so $x^{(i)}_{t,p}$ denotes the power of the p'th frequency bin in the audio frame at time t. $y^{(i)}$ is a transcription of the utterance $x^{(i)}$. One can see that $x^{(i)}$ is in fact a matrix of size depending on number of frequency bins and i'th audio clip length. We use the same number of power bins for all audio clips in dataset.

DS2 network's input is a time-series x and the output is a prediction over characters $p(l_t|x)$ for each output time-step. For English language possible values of l_t are:

- letters from 'a' to 'z';
- space;
- apostrophe;
- blank.

Adding non-letter characters allows to find word boundaries. Special symbol blank is outputted each time, when the network is unable to tell which character is most likely to occur for the current input spectrogram.

1.2. Layers

The model of the network is composed of one or more convolutional layers, followed by one or more bidirectional recurrent layers [4], followed by one or more fully connected layers. Activation function used throughout the network is the Clipped Rectified Linear (ReLU) function, given by the formula:

$$\sigma(x) = \begin{cases} 0 & \text{for } x < 0, \\ 20 & \text{for } x > 20, \\ x & \text{otherwise.} \end{cases}$$

The recurrent layers appear in a few different variants – standard recurrent layers or Long Short-Term Memory (LSTM) [5], or Gated Recurrent Units (GRU) [6]. After the recurrent layers and fully connected layers are applied, we count the output layer L as a softmax of the output of the last layer.

Softmax function $f : \mathbb{R}^k \to \mathbb{R}^k$ is defined by the formula:

$$f_i(v) = \frac{e^{v_i}}{\sum_{i=1}^{k} e^{v_j}},$$

and

$$f(v) = \big(f_1(v), f_2(v), \ldots, f_k(v)\big).$$

We basically apply exponential function to each outputed value, and then normalize these values to make sure, that probabilities sum up to 1. In our case $k = 29$, hence there is 29 possible output characters to distribute probability on (26 letters and 3 special symbols, as described in the previous section).

1.3. CTC Loss

To train a neural network we typically need a function that would tell us how good current network's output is. The lesser value this function has, the better results our model achieves. This kind of function is called a **loss function**. Usually, minimizing a value of the loss function is a main goal of the training.

Loss function used in DS2 is Connectionist Temporal Classification (CTC) [3]. To define this loss let us introduce an encoding of a text. Encoding of a given text S is done by replacing every character c in S by any number of characters c and blanks "-". Only restriction is that if there are two adjacent identical letters in S, they must be separated by a blank "-".

For instance, possible encodings of *"to"* are *"-tttooo"* and *"-tttoo-o"*, but only the latter could be an encoding of the word *"too"*.

Now we say that the probability of an actual transcription is the sum of probabilities of all possible encodings of the actual transcription, that have the same length as the output of the model. The loss function here is simply the negative logarithm of this probability. Having this we can count derivatives of this loss function with respect to model output and then apply *backpropagation through time* algorithm to train the network, which is just a standard backpropagation algorithm modified for recurrent neural networks.

1.4. Training

During the training data sets are divided into batches of some certain size. Batches are fed one by one into the model and after each one CTC loss is computed and backpropagated through the network, and then weights on the layers are updated. The speed in which changes are applied to the model depends on some chosen constant – **learning rate**. All data sets make up for an epoch of training and training lasts for many epochs – until loss doesn't decrease anymore or decreases very slightly. Accuracy of the model is measured by Word Error Rate (WER). It is a common metric of the performance of speech recognition systems. Here [8] many commercial ASR systems were compared using WER. To measure Word Error Rate, we compare the recognized word sequence with a reference word sequence by transforming the latter to the former and find:

- S – the number of substitutions,

- D – the number of deletions,
- I – the number of insertions,
- C – the number of correct words,
- N – the number of words in the reference ($N = S + D + C$).

Now the WER can be computed as:

$$WER = \frac{S + D + I}{N} = \frac{S + D + I}{S + D + C}$$

1.5. Generating transcription

In order to find Word Error Rate of the model, one has to generate some final transcription. It is generated based on both output of the model and a n-gram language model [9]. To generate the final transcription we search for transcription y that maximizes $Q(y)$, where Q is given as follows:

$$Q(y) = \log\left(\mathbb{P}_{ctc}(y|x)\right) + \alpha \cdot \log\left(\mathbb{P}_{lm}(y)\right) + \beta \cdot word\_count(y).$$

The term $\mathbb{P}_{ctc}(y|x)$ denotes the probability of y being a transcription of the utterance x as described in section 1.3 and the term $\mathbb{P}_{lm}(y)$ denotes the probability of the sequence y according to the language model. The weight α controls the relative contributions of the language model and the CTC network. The weight β encourages more words in the transcription. Both of those parameters are tunable. We use a beam search similar to one described in [7] to find the optimal transcription y.

Beam search is an algorithm that iterates over each time-step in the output of DS2 model, while remembering some constant number (let this number be BW – *Beam Width*) of the most probable transcriptions up to this point. When considering the next time-step, we create $BW \cdot C$ new transcriptions corresponding to appending every possible character (we denote number of them by C) to the end of each one of BW remembered transcriptions. After merging identical transcriptions and recalculating their probabilities we again keep BW best ones. In the end we have BW *most probable* transcriptions, so we just take the best one.

In order to efficiently merge transcriptions and recalculate new probabilities for them we must store probabilities $\mathbb{P}_b$, $\mathbb{P}_{nb}$, $\mathbb{P}_{lm}$ for all kept transcriptions, where:

- $\mathbb{P}_b(y, t)$ – the probability of the given transcription y at time-step t if the encoding ends with *blank*,
- $\mathbb{P}_{nb}(y, t)$ – the probability of the given transcription y at time-step t if the encoding doesn't end with *blank*,
- $\mathbb{P}_{lm}(y)$ – the probability of the given transcription y, according to the language model.

Nie podoba mi się używanie mathbb P na oznaczenie tych prawdopodobieństw, bo to sugeruje, że ono jest jakąś miarą, a jest tylko zwykłą funkcją.

1.6. Summary

To sum up workflow in DS2 is quite standard as for neural network. We prepare some data in the form described in 1.1. Next we forward propagate the data through the layers described in 1.2. After that we count CTC Loss (1.3.) and update model weights with specific learning rate (1.4.). To evaluate model we generate transcription from predictions given by the output layer using Beam search (1.5.) and count Word Error Rate metrics (1.4.). **TODO Hiperłącza**

Chapter 2

Experiments on architecture and hyperparameters

2.1. Regularization

2.1.1. Batch normalization

Batch normalization is a technique for improving the speed, performance, and stability of artificial neural networks, primarily introduced to reduce internal covariate shift [18]. Method bases on the idea of normalizing data within each mini-batch using stochastic optimizations for mean and variance.

Denote by B mini-batch of a training set, by m size of B. If input is d–dimensional: $x = (x^{(1)}, x^{(2)}, \ldots, x^{(d)})$, we calculate mean and variance as follows:

$$\mu^{(k)} = \frac{1}{m} \sum_{i=1}^{m} x_i^{(k)}$$

$$\sigma^{(k)^2} = \frac{1}{m} \sum_{i=1}^{m} \left(x_i^{(k)} - \mu^{(k)} \right)^2.$$

Next step is normalization:

$$\hat{x}_i^{(k)} = \frac{x_i^{(k)} - \mu^{(k)}}{\sqrt{\sigma^{(k)^2} + \epsilon}}$$

where ϵ is small constant added for numerical stability. Finally we scale and shift normalized data:

$$y_i^{(k)} = \gamma^{(k)} \hat{x}_i^{(k)} + \beta^{(k)}$$

where $\gamma^{(k)}$ and $\beta^{(k)}$ are learnable parameters.

2.1.2. Dropout

Dropout is regularization strategy patented by Google [20], it is used to remove neural network overfitting. The basic idea of this technique is multiplying input and hidden units, during training process, by randomly generated bit mask. [19]. It follows that some nodes will be inactive, set to 0, and model is prevented from complex co-adaptations.

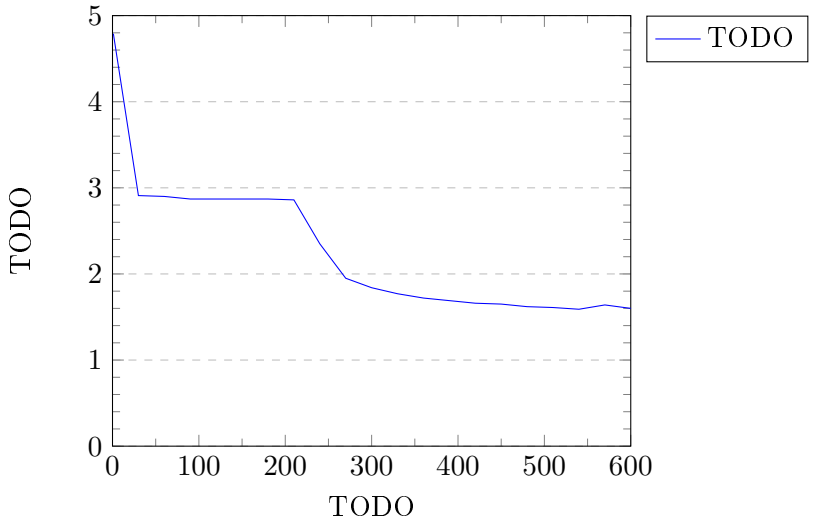

Figure 2.1: Best model performance in CTC Loss function with respect to training time in epochs.

2.1.3. L2 regularization

$L2$ regularization is a technique designed to prevent neural network from overfitting training dataset. The key idea behind this method is keeping network weights from growing too large unless it is really necessary. It can be realized by adding a term $\lambda||W||_F^2$ to the cost function that penalizes large weights, where $||W||_F$ is a Frobenius norm, and λ is parameter governing how strongly large weights are penalized [17].

2.2. Mixed precision training

In order to speed up training, we decided to try to reduce precision of floating-point calculations, where full precision is not really needed to achieve comparable accuracy.

Standard deep learning frameworks, such as PyTorch, use 32 bit floating points variables (FP32) to keep model parameters, and perform all operation. Massively parallel architecture of modern GPUs allows to, in the same time, perform much more arithmetic operations on 16 bit floating-point variables (FP16) than on 32 bit variables. Specifically, performance can be boosted by using Tensor Cores: hardware units prepared to accelerate matrix multiplications, introduced by Nvidia in Volta GPU architecture. [10] So reducing precision could not only reduce memory consumption but also measurably improve time needed for model convergence.

Common technique, used to improve speed, but not to break performance, is to keep model weights in two copies. One in FP32 and one in FP16. In each iteration we would use FP16 one to perform forward, backward propagation, and then to calculate gradients of every weight. Then we would use these gradients to update weights of original, FP32 model weights. Finally, we would make a new instance of FP16 model weight by cutting off precision of updated FP32 model. [11]

Another heuristic is to port to FP16 only operations, which performance would benefit most significantly (for instance transitions through fully connected and convolutional layers) and keep in FP32 operations, where precision might be an issue (softmax and batchnorm are here to mention). [11]

There is one more thing, we need to care about, while reducing precision. FP16 only 5 bits for exponent, and its representable range is quite narrow, so that in some cases gradients are small enough to be eventually rounded to 0. To prevent that one may multiply loss by a power of two, then perform backward propagation algorithm, and then scale back obtained gradients. This solution works pretty well in practice. However now we have to deal with another problem. Gradients calculated in the process may overflow. And even one overflow unables us to calculate gradients for parameters depended on that particular weight. So the scaling factor, rather than being one fixed power of two, need to be adjusted during training. Decreased after detection of overflow and increased if too many gradients are zeroed as a result of rounding.[13]

Speaking from more practical side, we used Apex AMP, which is a PyTorch Extension developed to ease mixed precision training. It supports automatic dynamic loss scaling, and both described above ways of performing mixed precision training.

Further description of this tool can be found in [12].

2.3. Language model and decoding predictions

N-gram language model is a model of a language in which we assign probabilities of occurrence to sequences of N words based on some large unlabelled text data. Incorporating model to prediction decoding helps avoid typos which are usually phonetically plausible and often occur on words that rarely or never occur in training set. We use ..... language model.

Beam size in beam search used in our model is ..... This size practically scales time of decoding linearly so we used such size that increasing it further wouldn't give any decrease in WER.

Performances with the language model in comparison to no language model are:

- ....% WER with LM; ...% WER without LM (model ....)

In the beginning we also experimented with an easier language model that assigned probabilities to sequences of letters as it was much smaller and simpler. Unfortunately WER was increasing after applying such LM so we abandoned this idea for the bigger and more reliable N-gram LM.

2.4. Recurrent unit type

In the section 1.2 we mentioned bidirectional recurrent layers. Three most commonly used ones are vanilla Recurrent Neural Networks (RNN) [14], Long Short Term Memory Unit (LSTM) [15] and Gated Recurrent Unit (GRU) [16]. We will deliver a brief description of these units based on [2].

Let's denote by l the layer number. A bidirectional recurrent layer h^l consist of a forward recurrent layer $\overrightarrow{h^l}$ and a backward recurrent layer $\overleftarrow{h^l}$. The forward and backward recurrent layer activations for time t are computed as $\overrightarrow{h_t^l} = g(\overrightarrow{h_{t-1}^l}, h_t^{l-1})$ and $\overleftarrow{h_t^l} = g(\overleftarrow{h_{t+1}^l}, h_t^{l-1})$. In the end, we add both partial layers to get $h^l = \overrightarrow{h^l} + \overleftarrow{h^l}$.

Function g mentioned before depends on specific recurrent unit type. For vanilla RNN it is just

$$\overrightarrow{h_t^l} = f(W^l h_t^{l-1} + \overrightarrow{U^l}\overrightarrow{h_{t-1}^l} + b^l)$$

and analogically for the backward recurrent layer. Here W^l and $\overrightarrow{U^l}$ are just weight matrices, b^l is a bias vector and f is an activation function (e.g. ReLU or tanh). For LSTM and GRU function g is much more complex to simulate some kind of *memory*.

This is a short summary of our experiments:

2.5. Performance and multi GPU scaling

For training, we used machine with Intel(R) Xeon(R) CPU E5-2698 v4 @ 2.20GHz, (which has 20 physical cores), and 256 GB of RAM. Is was equipped with four Nvidia Tesla V100 GPGPUs (32 GB of global memory each), connected by NVlink.

Figures below describe performance and scaling we achieved by utilizing more than one GPUs. Performance is measured by time in seconds it took for train for one epoch on full 1000h dataset.

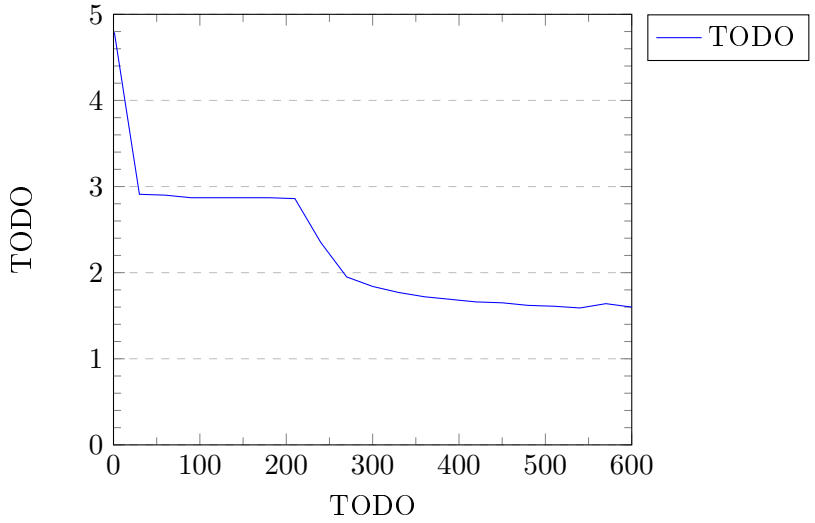

Figure 2.2: Best model performance in CTC Loss function with respect to training time in epochs.

And here is performance of mixed precision training.

2.6. Sortagrad and dataset sorting

2.7. Initialization

2.7.1. Xavier initialization

2.7.2. Random seed

2.8. Training dataset

2.9. Hyperparameters

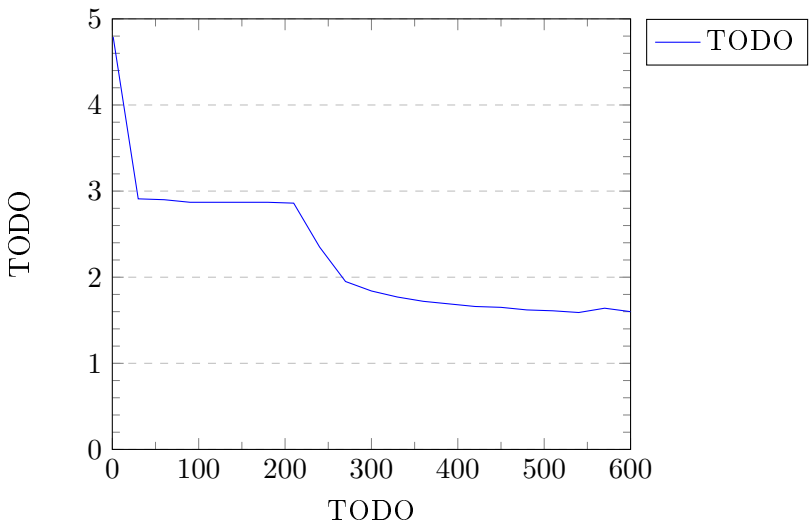

Figure 2.3: Best model performance in CTC Loss function with respect to training time in epochs.

Chapter 3

Conclusions

To sum up, we present `PyTorch` scripts for training DeepSpeech2 model for ASR. We also present already trained models for English and Polish as well as the results of our experiments justifying using specific hyperparameters and architecture solutions.

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