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Tutorial 10 - Self-Supervised, Representation & Transfer Learning



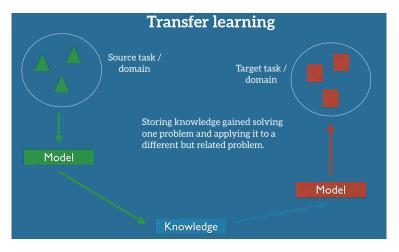
Agenda

- · Pre-trained Models and Transfer Learning
 - Transfer Learning
 - Pre-trained Models
- Representation Learning and Self-Supervised
 - Autoencoders
 - Contrastive Methods
- Recommended Videos
- Credits

```
In [27]: # imports for the tutorial
         import numpy as np
          import matplotlib.pyplot as plt
         import time
         import os
         import copy
          # pytorch imports
         import torch
         import torch.nn as nn
          import torch.nn.functional as F
         from torch.utils.data import DataLoader, Dataset
         from torchvision.datasets import ImageFolder
         from torchvision import models, transforms
         import torchvision
         # scikit-learn imports
         from sklearn.manifold import LocallyLinearEmbedding, Isomap, TSNE
         from sklearn.decomposition import PCA, KernelPCA
```

Transfer Learning

- Training deep neural networks has come a long way in past years, enabling learning good mappings from inputs to outputs, whether they are images, sentences, label predictions, etc. from large amounts of labelled data.
- · However, it is usually the case that our models struggle with generalization to unseen data during training.
- The traditional supervised learning paradigm breaks down when we don't have sufficient labeled data for the task or domain we want to train our model for.
 - For example, if we want to train a model to detect pedestrians on night-time images, we could, in theory, **apply a model that has been trained on a similar domain**, e.g. on day-time images. In practice, however, we often experience a deterioration or collapse in performance as the model has inherited the **bias** of its training data and does not know how to generalize to the new domain.
- Transfer learning allows us to deal with such scenarios by leveraging the already existing labelled data of some related task or domain.
 - We try to store this knowledge gained in solving the source task in the source domain and apply it to a new target task or domain.



· Image Source (https://ruder.io/transfer-learning/)

Transfer Learning Definition

- Transfer learning involves the concepts of a domain and a task.
- · For simplicity, we assume a binary classification setting.
- We denote the following: a domain $\mathcal D$ consists of a feature space $\mathcal X$ and a marginal probability distribution P(X) over the feature space, where $X=x_1,x_2,\ldots,x_n\in\mathcal X$.
- Given a domain, $\mathcal{D} = \{\mathcal{X}, P(X)\}$, a task \mathcal{T} consists of a label space \mathcal{Y} and a conditional probability distribution P(Y|X) that is typically learned from the training data consisting of pairs $x_i \in \mathcal{X}$ and $y_i \in \mathcal{Y}$.
 - In the binary classification case, $y_i \in \{True, False\}$.
- Given a source domain \mathcal{D}_S , a corresponding source task \mathcal{T}_S , as well as a target domain \mathcal{D}_T and a target task \mathcal{T}_T , the objective of transfer learning is to enable us to learn the **target conditional probability** distribution $P(Y_T|X_T)$ in \mathcal{D}_T with the information gained from \mathcal{D}_S and \mathcal{T}_S .
 - lacktriangledown We assume $\mathcal{D}_S
 eq \mathcal{D}_T$ ("domain adaption") or $\mathcal{T}_S
 eq \mathcal{T}_T$ ("transfer learning").
- In most cases, a limited number of labelled target examples, which is exponentially smaller than the number of labeled source examples are
 assumed to be available.



Transfer Learning Scenarios

- 1. $\mathcal{X}_S
 eq \mathcal{X}_T$ the feature spaces of the source and target domain are different.
 - Example: cross-lingual adaptation in a document classification task, the documents are written in two different languages.
- 2. $P(X_S) \neq P(X_T)$ the marginal probability distributions of source and target domain are different.
 - Example: domain adaptation in a document classification task, the documents (X) discuss different topics.
- 3. $\mathcal{Y}_S
 eq \mathcal{Y}_T$ the label spaces between the two tasks are different.
 - Example: documents need to be assigned different labels in the target task.
 - Usually happens with scenario 4 (as it is rare for two different tasks to have different label spaces, but exactly the same conditional
 probability distributions).
- 4. $P(Y_S|X_S)
 eq P(Y_T|X_T)$ the conditional probability distributions of the source and target tasks are different. Very common in practice.
 - Example: source and target documents are unbalanced with regard to their classes.



Transfer Learning Applications

- Sim2Real transferring from simulation to real environments. For many machine learning applications that rely on hardware for interaction, gathering data and training a model in the real world is either expensive, time-consuming, or simply too dangerous. It is thus advisable to gather data in some other, less risky way.
 - Learning from a simulation and applying the acquired knowledge to the real world is an instance of transfer learning scenario 2, as the feature spaces between source and target domain are the same (both generally rely on pixels), but the marginal probability distributions between simulation and reality are different, i.e. objects in the simulation and the source look different, although this difference diminishes as simulations get more realistic.
 - At the same time, the conditional probability distributions between simulation and real world might be different as the simulation is not able to fully replicate all reactions in the real world, e.g. a physics engine can not completely mimic the complex interactions of real-world objects.
 - Common applications include autonomus driving and robotics (where gathering data can be slow or dangerous).









- Image Source (https://arxiv.org/abs/1610.04286)
- Domain Adaptation transferring between domains that share some properties.
 - In vision, we commonly have many labelled data for some domain, but for the actual data that we care about, there are very few labels or none at all. Even if the training and the test data looks the same, the training data may still contain a bias that is imperceptible to humans but the model can still exploit it to overfit the training data.
 - In NLP, models trained on news data have difficulty coping with more novel text forms such as social media messages and the challenges they present. Even within one domain such as product reviews, people employ different words and phrases to express the same opinion. A model trained on one type of review should thus be able to disentangle the general and domain-specific opinion words that people use in order not to be confused by the shift in domain.

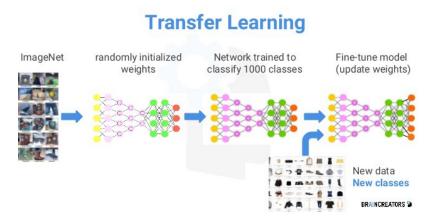


Transfer Learning with Pre-trained Models

- One of the fundamental requirements for transfer learning is the presence of models that perform well on source tasks.
- The two most common fields that build upon pre-trained models to transfer between tasks and domains are computer vision and NLP.

Using Pre-trained CNN Features

- Evidently, lower convolutional layers capture **low-level image features**, e.g. edges, while higher convolutional layers capture more complex details, such as body parts, faces, and other compositional features.
- The final fully-connected layers are generally assumed to capture information that is relevant for solving the respective task, e.g. classification.
- Representations that capture general information of how an image is composed and what combinations of edges and shapes it contains can
 be helpful in other tasks. This information is contained in one of the final convolutional layers or early fully-connected layers in large
 convolutional neural networks trained on ImageNet.
- For a new task, we can thus simply use the off-the-shelf features of a state-of-the-art CNN pre-trained on ImageNet and train a new model on these extracted features.
- In practice, we either keep the pre-trained parameters fixed or tune them with a small learning rate in order to ensure that we do not
 unlearn the previously acquired knowledge.



• Image Source (https://medium.com/datadriveninvestor/what-you-must-know-about-transfer-learning-4a6e4cb9fbad)



Transfer Learning Example with PyTorch

- We will follow examples by <u>Sasank Chilamkurthy (https://pytorch.org/tutorials/beginner/fransfer_learning_tutorial.html)</u> and <u>Nathan Inkawhich (https://pytorch.org/tutorials/beginner/finetuning_torchvision_models_tutorial.html)</u>.
- · We will train a classifier to distinguish between ants and bees.
 - The data can be downloaded from here: Download Link (https://download.pytorch.org/tutorial/hymenoptera_data.zip).
- · There are two major transfer learning scenarios:
 - Fine-tuning the ConvNet: Instead of random initialization, we initialize the network with a pretrained network, like VGG for example (which is trained on ImageNet 1000 dataset). Rest of the training looks as usual (except for usually using a lower learning rate).
 - ConvNet as fixed feature extractor: Here, we will freeze the weights for all of the network except that of the final fully connected layer. This last fully connected layer is replaced with a new one with random weights and only this layer is trained.

```
In [15]: # Data augmentation and normalization for training
                                      # Just normalization for validation
                                      {\tt data\_transforms} \; = \; \{ \;
                                                       'train': transforms.Compose([
                                                                     transforms.RandomResizedCrop(224).
                                                                     transforms.RandomHorizontalFlip(),
                                                                     transforms.ToTensor(),
                                                                     transforms.Normalize([0.485, 0.456, 0.406], [0.229, 0.224, 0.225])
                                                      1),
                                                         val': transforms.Compose([
                                                                     transforms.Resize(256),
                                                                     transforms.CenterCrop(224),
                                                                     transforms.ToTensor(),
                                                                     transforms.Normalize([0.485, 0.456, 0.406], [0.229, 0.224, 0.225])
                                                      ]),
                                      }
                                      batch_size = 4
                                      data_dir = './datasets/hymenoptera_data'
                                      image\_datasets = \{x: \; ImageFolder(os.path.join(data\_dir, \; x), \; data\_transforms[x]) \; \\ \textit{for} \; x \; \\ \textit{in} \; ['train', \; 'val']\} \\ image\_datasets = \{x: \; ImageFolder(os.path.join(data\_dir, \; x), \; data\_transforms[x]) \; \\ \textit{for} \; x \; \\ \textit{in} \; ['train', \; 'val']\} \\ image\_datasets = \{x: \; ImageFolder(os.path.join(data\_dir, \; x), \; data\_transforms[x]) \; \\ \textit{for} \; x \; \\ \textit{in} \; ['train', \; 'val']\} \\ image\_datasets = \{x: \; ImageFolder(os.path.join(data\_dir, \; x), \; data\_transforms[x]) \; \\ \textit{for} \; x \; \\ \textit{in} \; ['train', \; 'val']\} \\ image\_datasets = \{x: \; ImageFolder(os.path.join(data\_dir, \; x), \; data\_transforms[x]) \; \\ \textit{for} \; x \; \\ \textit{in} \; ['train', \; 'val']\} \\ image\_datasets = \{x: \; ImageFolder(os.path.join(data\_dir, \; x), \; data\_transforms[x]) \; \\ \textit{for} \; x \; \\ \textit{fo
                                      dataloaders = {x: DataLoader(image_datasets[x], batch_size=batch_size,
                                                                                                                                                          shuffle=True, num_workers=4) for x in ['train', 'val']}
                                      dataset_sizes = {x: len(image_datasets[x]) for x in ['train', 'val']}
                                      class_names = image_datasets['train'].classes
                                      device = torch.device("cuda:0" if torch.cuda.is_available() else "cpu")
                                     print(device)
```

cuda:0

```
In [10]: def imshow(inp, title=None):
    """Imshow for Tensor."""
    inp = inp.numpy().transpose((1, 2, 0))
    mean = np.array([0.485, 0.456, 0.406])
    std = np.array([0.229, 0.224, 0.225])
    inp = std * inp + mean
    inp = np.clip(inp, 0, 1)
    fig = plt.figure(figsize=(5, 8))
    ax = fig.add_subplot(111)
    ax.imshow(inp)
    if title is not None:
        ax.set_title(title)
    ax.set_axis_off()
```

```
In [11]: # Let's visualize a few training images so as to understand the data augmentations.
# Get a batch of training data
inputs, classes = next(iter(dataloaders['train']))

# Make a grid from batch
out = torchvision.utils.make_grid(inputs)
imshow(out, title=[class_names[x] for x in classes])
```

['ants', 'bees', 'bees']

Set Model Parameters' .requires_grad attribute

- The following helper function sets the .requires_grad attribute of the parameters in the model to False when we are feature extracting.
- By default, when we load a pretrained model all of the parameters have .requires_grad=True, which is fine if we are training from scratch or finetuning
- However, if we are feature extracting and only want to compute gradients for the newly initialized layer then we want all of the other parameters to not require gradients.

```
In [13]: def set_parameter_requires_grad(model, feature_extracting=False):
    if feature_extracting:
        for param in model.parameters():
            param.requires_grad = False
    else:
        for param in model.parameters():
            param.requires_grad = True
```

Initialize and Reshape the Networks

- Recall, the final layer of a CNN model, which is often an FC layer, has the same number of nodes as the number of output classes in the dataset.
- · Since all of the following models have been pretrained on ImageNet, they all have output layers of size 1000, one node for each class.
- The goal here is to reshape the last layer to have the same number of inputs as before, AND to have the same number of outputs as the number of classes in the dataset.
- When feature extracting, we only want to update the parameters of the last layer, or in other words, we only want to update the parameters for the layer(s) we are reshaping.
- Therefore, we do not need to compute the gradients of the parameters that we are not changing, so for efficiency we set the <code>.requires_grad</code> attribute to <code>False</code>.
- This is important because by default, this attribute is set to True. Then, when we initialize the new layer and by default the new parameters have .requires grad=True so only the new layer's parameters will be updated.
- When we are fine-tuning we can leave all of the .required_grad 's set to the default of True .

```
In [20]: | def initialize_model(model_name, num_classes, feature_extract, use_pretrained=True):
             # Initialize these variables which will be set in this if statement. Each of these
             # variables is model specific.
             model_ft = None
             input_size = 0  # image size, e.g. (3, 224, 224)
             if model name == "resnet":
                  """ Resnet18
                 model_ft = models.resnet18(pretrained=use_pretrained)
                 set_parameter_requires_grad(model_ft, feature_extract)
                 num_ftrs = model_ft.fc.in_features
                 model_ft.fc = nn.Linear(num_ftrs, num_classes) # replace the last FC layer
                 input_size = 224
             elif model_name == "alexnet":
                 """ Alexnet
                 model_ft = models.alexnet(pretrained=use_pretrained)
                 set_parameter_requires_grad(model_ft, feature_extract)
                 num_ftrs = model_ft.classifier[6].in_features
                 model_ft.classifier[6] = nn.Linear(num_ftrs, num_classes)
                 input_size = 224
             elif model_name == "vgg":
                 """ VGG16
                 model_ft = models.vgg16(pretrained=use_pretrained)
                 set_parameter_requires_grad(model_ft, feature_extract)
                 num_ftrs = model_ft.classifier[6].in_features
                 model_ft.classifier[6] = nn.Linear(num_ftrs, num_classes)
                 input\_size = 224
             elif model_name == "squeezenet":
                  """ Squeezenet
                 model_ft = models.squeezenet1_0(pretrained=use_pretrained)
                 set_parameter_requires_grad(model_ft, feature_extract)
                 model_ft.classifier[1] = nn.Conv2d(512, num_classes, kernel_size=(1,1), stride=(1,1))
                 model_ft.num_classes = num_classes
                 input size = 224
             elif model_name == "densenet":
                 """ Densenet
                 model_ft = models.densenet121(pretrained=use_pretrained)
                 set_parameter_requires_grad(model_ft, feature_extract)
                 num_ftrs = model_ft.classifier.in_features
                 model_ft.classifier = nn.Linear(num_ftrs, num_classes)
                 input_size = 224
             else:
                 raise NotImplementedError
             return model_ft, input_size
```

```
In [21]: # Models to choose from [resnet, alexnet, vgg, squeezenet, densenet]
model_name = "vgg"

# Number of classes in the dataset
num_classes = 2

# Batch size for training (change depending on how much memory you have)
batch_size = 8

# Number of epochs to train for
num_epochs = 15

# Flag for feature extracting. When False, we fine-tune the whole model,
# when True we only update the reshaped layer params
feature_extract = True
```

```
In [22]: # Initialize the model for this run
         model ft, input size = initialize model(model name, num classes, feature extract, use pretrained=True)
         # Print the model we just instantiated
         print(model ft)
         VGG(
           (features): Sequential(
             (0): Conv2d(3, 64, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (1): ReLU(inplace=True)
             (2): Conv2d(64, 64, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (3): ReLU(inplace=True)
             (4): MaxPool2d(kernel_size=2, stride=2, padding=0, dilation=1, ceil_mode=False)
             (5): Conv2d(64, 128, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (6): ReLU(inplace=True)
             (7): Conv2d(128, 128, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (8): ReLU(inplace=True)
             (9): MaxPool2d(kernel_size=2, stride=2, padding=0, dilation=1, ceil_mode=False)
             (10): Conv2d(128, 256, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (11): ReLU(inplace=True)
             (12): Conv2d(256, 256, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (13): ReLU(inplace=True)
             (14): Conv2d(256, 256, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (15): ReLU(inplace=True)
             (16): MaxPool2d(kernel_size=2, stride=2, padding=0, dilation=1, ceil_mode=False)
             (17): Conv2d(256, 512, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (18): ReLU(inplace=True)
             (19): Conv2d(512, 512, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (20): ReLU(inplace=True)
             (21): Conv2d(512, 512, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (22): ReLU(inplace=True)
             (23): MaxPool2d(kernel_size=2, stride=2, padding=0, dilation=1, ceil_mode=False)
             (24): Conv2d(512, 512, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (25): ReLU(inplace=True)
             (26): Conv2d(512, 512, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (27): ReLU(inplace=True)
             (28): Conv2d(512, 512, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
             (29): ReLU(inplace=True)
             (30): MaxPool2d(kernel_size=2, stride=2, padding=0, dilation=1, ceil_mode=False)
           (avgpool): AdaptiveAvgPool2d(output_size=(7, 7))
           (classifier): Sequential(
             (0): Linear(in_features=25088, out_features=4096, bias=True)
             (1): ReLU(inplace=True)
             (2): Dropout(p=0.5, inplace=False)
             (3): Linear(in_features=4096, out_features=4096, bias=True)
             (4): ReLU(inplace=True)
             (5): Dropout(p=0.5, inplace=False)
             (6): Linear(in_features=4096, out_features=2, bias=True)
In [24]: model_ft = model_ft.to(device)
         # Gather the parameters to be optimized/updated in this run. If we are
         # fine-tuning we will be updating all parameters. However, if we are
         # doing feature extract method, we will only update the parameters
         # that we have just initialized, i.e. the parameters with requires_grad
         # is True.
         params_to_update = model_ft.parameters()
         print("Params to learn:")
         if feature_extract:
             params_to_update = []
             for name,param in model_ft.named_parameters():
                 if param.requires_grad == True:
                     params_to_update.append(param)
                     print("\t",name)
         else:
             for name,param in model_ft.named_parameters():
                 if param.requires_grad == True:
                     print("\t",name)
         # Observe that all parameters are being optimized
         optimizer_ft = torch.optim.SGD(params_to_update, lr=0.001, momentum=0.9)
         Params to learn:
```

classifier.6.weight classifier.6.bias

```
....
In [28]:
         Training function
         def train_model(model, dataloaders, criterion, optimizer, num_epochs=25):
             since = time.time()
             val_acc_history = []
             best_model_wts = copy.deepcopy(model.state_dict())
             best acc = 0.0
             for epoch in range(num_epochs):
                 print('Epoch {}/{}'.format(epoch, num_epochs - 1))
                 print('-' * 10)
                  # Each epoch has a training and validation phase
                 for phase in ['train', 'val']:
                     if phase == 'train':
                         model.train() # Set model to training mode
                      else:
                         model.eval() # Set model to evaluate mode
                     running_loss = 0.0
                     running_corrects = 0
                      # Iterate over data.
                     for inputs, labels in dataloaders[phase]:
                         inputs = inputs.to(device)
                         labels = labels.to(device)
                         # forward
                          # track history if only in train
                         with torch.set_grad_enabled(phase == 'train'):
                             # Get model outputs and calculate loss
                             outputs = model(inputs)
                             loss = criterion(outputs, labels)
                              _, preds = torch.max(outputs, 1)
                              # backward + optimize only if in training phase
                              if phase == 'train':
                                  # zero the parameter gradients
                                  optimizer.zero_grad()
                                 loss.backward()
                                  optimizer.step()
                         # statistics
                         running_loss += loss.item() * inputs.size(0)
                         running_corrects += torch.sum(preds == labels.data)
                     epoch_loss = running_loss / len(dataloaders[phase].dataset)
                     epoch_acc = running_corrects.double() / len(dataloaders[phase].dataset)
                     print('{} Loss: {:.4f} Acc: {:.4f}'.format(phase, epoch_loss, epoch_acc))
                      # deep copy the model
                     if phase == 'val' and epoch_acc > best_acc:
                         best_acc = epoch_acc
                         best_model_wts = copy.deepcopy(model.state_dict())
                     if phase == 'val':
                         val_acc_history.append(epoch_acc)
                 print()
             time_elapsed = time.time() - since
             print('Training complete in {:.0f}m {:.0f}s'.format(time_elapsed // 60, time_elapsed % 60))
             print('Best val Acc: {:4f}'.format(best_acc))
             # Load best model weights
             model.load_state_dict(best_model_wts)
             return model, val_acc_history
```

Epoch 0/14

train Loss: 0.2460 Acc: 0.9057 val Loss: 0.1525 Acc: 0.9542

Epoch 1/14

train Loss: 0.2114 Acc: 0.9139 val Loss: 0.1093 Acc: 0.9477

Epoch 2/14

train Loss: 0.1623 Acc: 0.9221 val Loss: 0.1752 Acc: 0.9412

Epoch 3/14

train Loss: 0.2186 Acc: 0.9262 val Loss: 0.1530 Acc: 0.9477

Epoch 4/14

train Loss: 0.1423 Acc: 0.9467 val Loss: 0.1667 Acc: 0.9412

Epoch 5/14

train Loss: 0.1301 Acc: 0.9508 val Loss: 0.1108 Acc: 0.9673

Epoch 6/14

train Loss: 0.1602 Acc: 0.9303 val Loss: 0.1001 Acc: 0.9412

Epoch 7/14

train Loss: 0.1888 Acc: 0.9385 val Loss: 0.1254 Acc: 0.9542

Epoch 8/14

train Loss: 0.1432 Acc: 0.9508 val Loss: 0.1194 Acc: 0.9608

Epoch 9/14

train Loss: 0.1139 Acc: 0.9631 val Loss: 0.1230 Acc: 0.9608

Epoch 10/14

train Loss: 0.2273 Acc: 0.9467 val Loss: 0.1154 Acc: 0.9608

Epoch 11/14

train Loss: 0.1507 Acc: 0.9385 val Loss: 0.1096 Acc: 0.9608

Epoch 12/14

train Loss: 0.0876 Acc: 0.9590 val Loss: 0.1767 Acc: 0.9608

Epoch 13/14

train Loss: 0.3315 Acc: 0.9344 val Loss: 0.1371 Acc: 0.9608

Epoch 14/14

train Loss: 0.0936 Acc: 0.9672 val Loss: 0.1540 Acc: 0.9673

Training complete in 3m 39s Best val Acc: 0.967320

AX

Pre-training for Natural Language Processing

- One of the biggest challenges in NLP is the shortage of training data.
- Because NLP is a diversified field with many distinct tasks, most task-specific datasets contain only a few thousand or a few hundred thousand human-labeled training examples.
- As big companies like Google and OpenAl have shown, modern deep learning-based NLP models see benefits from much larger amounts of data, improving when trained on millions, or billions, of annotated training examples.
- Large pre-trained models on the enormous amount of unannotated text on the web can then be fine-tuned on small-data NLP tasks like
 question answering and sentiment analysis, resulting in substantial accuracy improvements compared to training on these datasets from
 scratch.
- Bidirectional Encoder Representations from Transformers (BERT), Google a Transformer-based machine learning technique for natural language processing (NLP) pre-training developed by Google. The idea is to mask certain words and then try to predict them. The original English-language BERT model comes with two pre-trained general types:
 - (1) the BERT_{BASE} model, a 12-layer, 768-hidden, 12-heads, 110M parameter neural network architecture.
 - (2) the $BERT_{LARGE}$ model, a 24-layer, 1024-hidden, 16-heads, 340M parameter neural network architecture.
 - Both of which were trained on the BooksCorpus dataset with 800M words, and a version of the English Wikipedia with 2,500M words.
- BERT uses the straightforward technique of masking out some of the words in the input and then condition each word bidirectionally to predict the masked words and also learns to model relationships between sentences by pre-training on a very simple task that can be generated from any text corpus: Given two sentences A and B, is B the actual next sentence that comes after A in the corpus, or just a random sentence?

```
{\bf Input}: The man went to the {\rm [MASK]}_1 . He bought a {\rm [MASK]}_2 of milk . Labels: {\rm [MASK]}_1 = store; {\rm [MASK]}_2 = gallon
```

```
Sentence A = The man went to the store.

Sentence B = He bought a gallon of milk.

Label = IsNextSentence
```

Sentence A = The man went to the store.

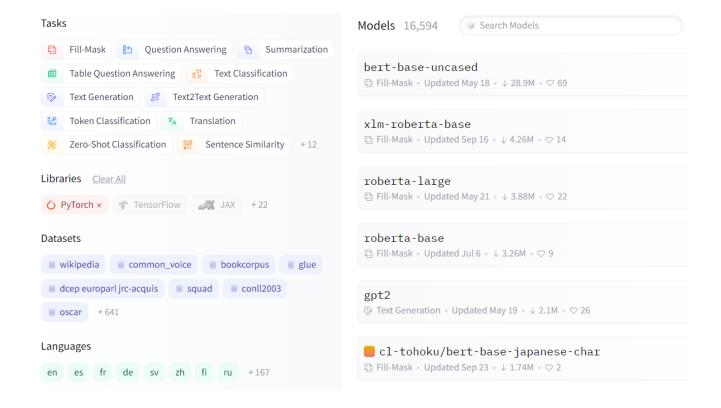
Sentence B = Penguins are flightless.

Label = NotNextSentence

• Image Source (https://ai.googleblog.com/2018/11/open-sourcing-bert-state-of-art-pre.html)

Pre-trained Models for NLP in PyTorch

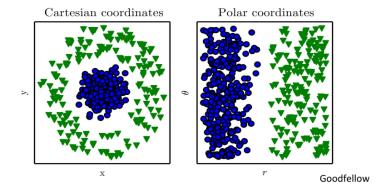
- <u>HuggingFace (https://huggingface.co/)</u> is a copmpany that is dedicated to publishing all of the available pretrained models and it works in PyTorch as well <u>HuggingFace Transformers (https://github.com/huggingface/transformers)</u>
- Examples with PyTorch (https://pytorch.org/hub/huggingface_pytorch-transformers/)





Representation and Self-Supervised Learning

- How do we learn rich and useful features from raw unlabeled data that can be useful for several downstream tasks (e.g., classification, reinforcement learning...)?
- What are the various general tasks that can be used to learn representations from unlabeled data?
- · The way we represent the data has a great impact on the performance and compelxity.
- Such representations can be learned in typical unsupervised settings, or in a self-supervised manner.

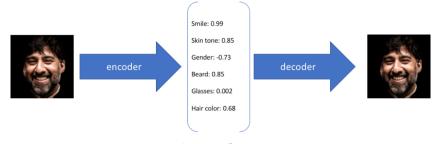


- **Deep Unsupervised Learning** learn representations without lables, subset of deep learning, which is a subset of representation learning, which is a subset of machine learning.
- Self-supervised Learning often used interchangeably with unsupervised learning. Self-supervised: create your own supervision through pretext tasks.

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Deep Unsupervised Learning - Deep Autoencoders

- Most of the natural data is high-dimensional, such as images. Consider the MNIST (hand-written digits) dataset, where each image has 28x28 = 784 pixels, which means it can be represented by a vector of length 784.
 - But do we really need 784 values to represent a digit? The answer is probably no. We believe that the data lies on a low-dimensional space which is enough to describe the observasions. In the case of MNIST, we can choose to represent digits as one-hot vectors, which means we only need 10 dimensions. So we can **encode** high-dimensional observations in a low-dimensional space.
 - But how can we learn meaningful low-deminsional representations? The general idea is to reconstruct or, decode the low-dimensional representation to the high-dimensional representation, and use the reconstruction error to find the best representations (using the gradients of the error). This is the core idea behind autoencoders.
 - Autoencoders models which take data as input and discover some latent state representation of that data. The input data is converted into an encoding vector where each dimension represents some learned attribute about the data. The most important detail to grasp here is that our encoder network is outputting a single value for each encoding dimension. The decoder network then subsequently takes these values and attempts to recreate the original input. Autoencoders have three parts: an encoder, a decoder, and a 'loss' function that maps one to the other. For the simplest autoencoders the sort that compress and then reconstruct the original inputs from the compressed representation we can think of the 'loss' as describing the amount of information lost in the process of reconstruction.
 - · Illustration:



Latent attributes

· The basic architecture of an autoencoder:

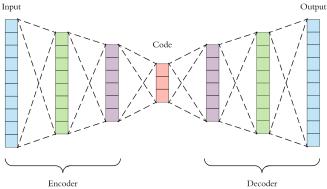


 Image from Applied Deep Learning (https://towardsdatascience.com/applied-deep-learning-part-3-autoencoders-1c083af4d798) by Arden Dertat

Let's implement it in PyTorch using what we have learnt so far!

```
In [14]: class AutoEncoder(nn.Module):
             def __init__(self, input_dim=28*28, hidden_dim=256, latent_dim=10):
                 super(AutoEncoder, self).__init__()
                 self.input_dim = input_dim
                 self.hidden dim = hidden dim
                 self.latent_dim = latent_dim
                 # define the encoder
                 self.encoder = nn.Sequential(nn.Linear(self.input_dim, self.hidden_dim),
                                               nn.ReLU(), nn.Linear(self.hidden_dim, self.hidden_dim),
                                               nn.ReLU(),
                                               nn.Linear(self.hidden_dim, self.latent_dim)
                 # define decoder
                 self.decoder = nn.Sequential(nn.Linear(self.latent_dim, self.hidden_dim),
                                               nn.ReLU(),
                                               nn.Linear(self.hidden_dim, self.hidden_dim),
                                               nn.ReLU(),
                                               nn.Linear(self.hidden_dim, self.input_dim),
                                               nn.Sigmoid())
             def forward(self,x):
                 x = self.encoder(x)
                 x = self.decoder(x)
                 \textbf{return} \ x
             def get_latent_rep(self, x):
                 return self.encoder(x)
```

```
In [15]: # hyper-parameters:
    num_epochs = 5
    learning_rate = 0.001

# Device configuration, as before
    device = torch.device('cuda:0' if torch.cuda.is_available() else 'cpu')

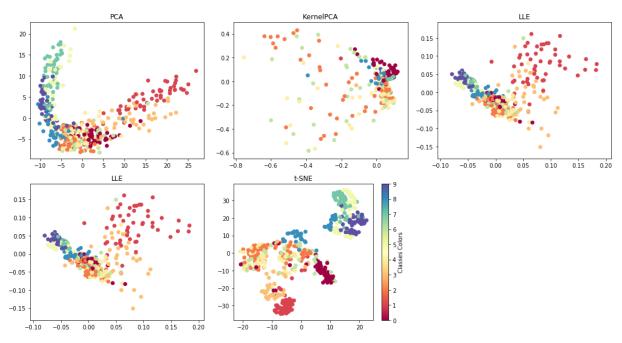
# create model, send it to device
    model = AutoEncoder(input_dim=28*28, hidden_dim=128, latent_dim=10).to(device)

# Loss and optimizer
    criterion = nn.BCELoss() # binary cross entropy, as pixels are in [0,1]
    optimizer = torch.optim.Adam(model.parameters(), lr=learning_rate)
```

```
In [16]: # Train the model
          total_step = len(fmnist_train_loader)
          for epoch in range(num_epochs):
              for i, (images, labels) in enumerate(fmnist_train_loader):
                   # each i is a batch of 128 samples
                   images = images.to(device).view(batch_size, -1)
                   # Forward pass
                   outputs = model(images)
                   loss = criterion(outputs, images)
                   # Backward and optimize - ALWAYS IN THIS ORDER!
                   optimizer.zero_grad()
                   loss.backward()
                   optimizer.step()
                   if (i + 1) % 100 == 0:
                       print ('Epoch [{}/{}], Step [{}/{}], Loss: {:.4f}'
                               .format(epoch + 1, num_epochs, i + 1, total_step, loss.item()))
          Epoch [1/5], Step [100/468], Loss: 0.3832
Epoch [1/5], Step [200/468], Loss: 0.3317
          Epoch [1/5], Step [300/468], Loss: 0.3360
          Epoch [1/5], Step [400/468], Loss: 0.3037
          Epoch [2/5], Step [100/468], Loss: 0.3093
          Epoch [2/5], Step [200/468], Loss: 0.2943
          Epoch [2/5], Step [300/468], Loss: 0.3223
          Epoch [2/5], Step [400/468], Loss: 0.3077
          Epoch [3/5], Step [100/468], Loss: 0.2953
          Epoch [3/5], Step [200/468], Loss: 0.3018
          Epoch [3/5], Step [300/468], Loss: 0.3083
Epoch [3/5], Step [400/468], Loss: 0.3133
          Epoch [4/5], Step [100/468], Loss: 0.3082
          Epoch [4/5], Step [200/468], Loss: 0.3007
          Epoch [4/5], Step [300/468], Loss: 0.2898
          Epoch [4/5], Step [400/468], Loss: 0.2976
          Epoch [5/5], Step [100/468], Loss: 0.3038
          Epoch [5/5], Step [200/468], Loss: 0.2896
Epoch [5/5], Step [300/468], Loss: 0.2972
          Epoch [5/5], Step [400/468], Loss: 0.2990
In [17]: # let's see some of the reconstructions
          model.eval() # put in evaluation mode - no gradients
          examples = enumerate(fmnist_test_loader)
          batch_idx, (example_data, example_targets) = next(examples)
          print("shape: \n", example_data.shape)
          fig = plt.figure()
          for i in range(3):
              ax = fig.add_subplot(2,3,i+1)
              ax.imshow(example\_data[i][\emptyset], \ cmap='gray', \ interpolation='none')
              ax.set_title("Ground Truth: {}".format(example_targets[i]))
              ax.set_axis_off()
              ax = fig.add_subplot(2,3,i+4)
              recon_img = model(example_data[i][0].view(1, -1).to(device)).data.cpu().numpy().reshape(28, 28)
              ax.imshow(recon_img, cmap='gray')
              ax.set_title("Reconstruction of: {}".format(example_targets[i]))
              ax.set_axis_off()
          plt.tight_layout()
          shape:
           torch.Size([128, 1, 28, 28])
              Ground Truth: 9
                                  Ground Truth: 2
                                                     Ground Truth: 1
            Reconstruction of: 9
                                Reconstruction of: 2
                                                   Reconstruction of: 1
```

```
In [19]: | fig = plt.figure(figsize=(15,8))
          # PCA
          t0 = time.time()
          x_pca = PCA(n_components).fit_transform(latent_X)
         t1 = time.time()
         print("PCA time: %.2g sec" % (t1 - t0))
          ax = fig.add_subplot(2, 3, 1)
         ax.scatter(x_pca[:, 0], x_pca[:, 1], c=labels, cmap=plt.cm.Spectral)
ax.set_title('PCA')
          # KPCA
         t0 = time.time()
          x_kpca = KernelPCA(n_components, kernel='rbf').fit_transform(latent_X)
          t1 = time.time()
          print("KPCA time: %.2g sec" % (t1 - t0))
          ax = fig.add_subplot(2, 3, 2)
          ax.scatter(x\_kpca[:, \ 0], \ x\_kpca[:, \ 1], \ c=labels, \ cmap=plt.cm.Spectral)
          ax.set_title('KernelPCA')
          # LLE
         t0 = time.time()
          x_lle = LocallyLinearEmbedding(n_neighbors, n_components, eigen_solver='auto').fit_transform(latent_X)
          t1 = time.time()
          print("LLE time: %.2g sec" % (t1 - t0))
          ax = fig.add_subplot(2, 3, 3)
          ax.scatter(x_lle[:, 0], x_lle[:, 1], c=labels, cmap=plt.cm.Spectral)
          ax.set_title('LLE')
          # Isomap
         t0 = time.time()
          x_isomap = Isomap(n_neighbors, n_components).fit_transform(latent_X)
          t1 = time.time()
         print("Isomap time: %.2g sec" % (t1 - t0))
          ax = fig.add_subplot(2, 3, 4)
         ax.scatter(x_lle[:, 0], x_lle[:, 1], c=labels, cmap=plt.cm.Spectral)
ax.set_title('LLE')
          # t-SNE
         t0 = time.time()
         x_tsne = TSNE(n_components).fit_transform(latent_X)
         t1 = time.time()
          print("t-SNE time: %.2g sec" % (t1 - t0))
         ax = fig.add_subplot(2, 3, 5)
          scatter = ax.scatter(x_tsne[:, 0], x_tsne[:, 1], c=labels, cmap=plt.cm.Spectral)
          ax.set_title('t-SNE')
          bounds = np.linspace(0, 10, 11)
          cb = plt.colorbar(scatter, spacing='proportional', ticks=bounds)
          cb.set_label('Classes Colors')
         plt.tight_layout()
```

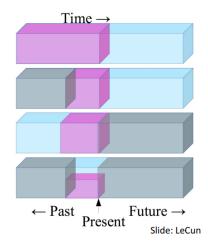
PCA time: 0.001 sec KPCA time: 0.013 sec LLE time: 0.077 sec Isomap time: 0.12 sec t-SNE time: 2.7 sec





Self-Supervised Learning

- A version of unsupervised learning where data provides the supervision.
- · Idea: withhold some part of the data and then task a neural network to predict it from the remaining parts.
- Details decide what proxy loss or pretext task the network tries to solve, and depending on the quality of the task, good semantic features can be obtained without actual labels.
- Advantages over supervised learning:
 - Large cost of producing a new dataset for each task (prepare labeling manuals, categories, hiring humans, creating GUIs, storage pipelines, etc).
 - Good supervision may not be cheap (e.g., medicine, legal).
 - Take advantage of vast amount of unlabeled data on the internet (images, videos, language).
 - Predict any part of the input from any other part.
 - ► Predict the future from the past.
 - Predict the future from the recent past.
 - ► Predict the past from the present.
 - ► Predict the top from the bottom.
 - Predict the occluded from the visible
 - Pretend there is a part of the input you don't know and predict that.



Self-Supervised Learning Methods

- · Reconstruct from a corrupted (or partial) version
 - Denoising Autoencoders
 - In-painting
 - Colorization, Split-Brain Autoencoder
- · Visual common sense tasks
 - Relative patch prediction
 - Jigsaw puzzles
 - Rotation prediction
- Contrastive Learning (our focus)
 - word2vec
 - Contrastive Predictive Coding (CPC)
 - Instance Discrimination
 - Simple Framework for Contrastive Learning of Visual Representations (SimCLR), Momentum Contrast (MoCo), Bootstrap Your Own Latent (BYOL)

Code Demos - Self-Supervised Learning Demos

(https://colab.research.google.com/github/rll/deepul/blob/master/demos/lecture7_selfsupervised_demos.ipynb)

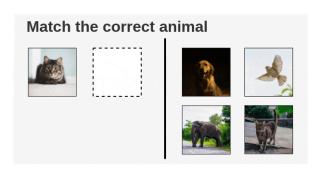


- Contrastive learning is an approach to formulate the task of finding similar and dissimilar things for a ML model (basically what classification does when given labels).
- · Contrastive methods, as the name implies, learn representations by contrasting positive and negative examples.
- · Using this approach, one can train a machine learning model to classify between similar and dissimilar images.
- More formally, for any data point x, contrastive methods aim to learn an encoder f such that:
 - x^+ is a data point similar to x, referred to as a *positive* sample.
 - x^- is a data point dissimilar to x, referred to as a *negative* sample.
 - The **score function** is a metric that measures the similarity between two features:

$$score(f(x),f(x^+)) >> score(f(x),f(x^-))$$



Loss measured in the representation space Examples: TCN, CPC, Deep-InfoMax



• Image Source (https://analyticsindiamag.com/contrastive-learning-self-supervised-ml)

· The most common loss function to implement the score paradigm is InfoNCE loss, which looks similar to softmax.

$$\mathcal{L}_{N} = -\mathbb{E}_{X} \left[\log \frac{\exp \left(f(x)^{T} f\left(x^{+}\right) \right)}{\exp \left(f(x)^{T} f\left(x^{+}\right) \right) + \sum_{j=1}^{N-1} \exp \left(f(x)^{T} f\left(x_{j}\right) \right)} \right]$$

• The denominator terms consist of one positive sample, and N-1 negative samples.

Contrastive Predictive Coding (CPC)

- **Contrastive Predictive Coding (CPC)** (https://arxiv.org/abs/1807.03748) learns self-supervised representations by **predicting the future** in a learned *latent space* by using powerful autoregressive models.
- The model uses a probabilistic contrastive loss which induces the latent space to capture information that is maximally useful to predict future samples.

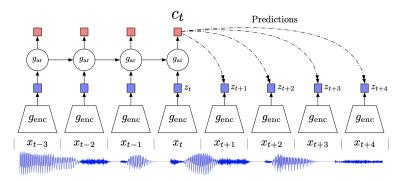


Figure 1: Overview of Contrastive Predictive Coding, the proposed representation learning approach. Although this figure shows audio as input, we use the same setup for images, text and reinforcement learning.

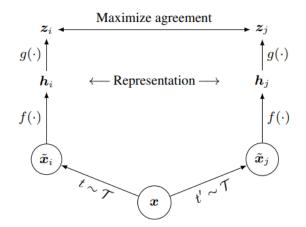
- First, a non-linear encoder g_{enc} maps the input sequence of observations x_t to a sequence of latent representations $z_t = g_{enc}(x_t)$, potentially with a lower temporal resolution.
- Next, an autoregressive model g_{ar} summarizes all $z \le t$ in the latent space and produces a context latent representation $c_t = g_{ar}(z \le t)$.
- A density ratio f is modelled which preserves the mutual information between x_{t+k} and c_t as follows:

$$f_k(x_{t+k}, c_t) = \exp(z_{t+k}^T W_k c_t) \propto rac{p(x_{t+k}|c_t)}{p(x_{t+k})}$$

- lacksquare Note that the density ratio f can be unnormalized (does not have to integrate to 1).
- W_k are learned weights.
- · Any type of encoder and autoregressive can be used.
 - For example: strided convolutional layers with residual blocks and GRUs.
- The encoder and autoregressive models are trained to minimize an InfoNCE loss.
- PyTorch Code (https://github.com/jefflai108/Contrastive-Predictive-Coding-PyTorch)

Simple Framework for Contrastive Learning of Visual Representations (SimCLR)

- **Simple Framework for Contrastive Learning of Visual Representations (SimCLR)** (https://arxiv.org/abs/2002.05709) is a framework for contrastive learning of visual representations.
- It learns representations by maximizing agreement between differently augmented views of the same data example via a contrastive loss in the latent space.



- A stochastic data augmentation module that transforms any given data example randomly resulting in two correlated views of the same example, denoted \tilde{x}_i and \tilde{x}_j , which is considered a **positive pair**.
- SimCLR sequentially applies three simple augmentations: random cropping followed by resize back to the original size, random color distortions, and random Gaussian blur. The authors find **random crop and color distortion** is crucial to achieve good performance.
- A neural network base encoder $f(\cdot)$ that extracts **representation vectors** from augmented data examples. The framework allows various choices of the network architecture without any constraints.
 - ullet For simplicity ResNet is used to obtain $h_i=f(ilde x_i)\in\mathcal{R}^d$ where h_i is the output after the average pooling layer.
- A small neural network projection head $g(\cdot)$ that maps representations to the space where contrastive loss is applied.
- MLP with one hidden layer is used to obtain $z_i = g(h_i)$.
- The authors find it beneficial to define the contrastive loss on z_i 's rather than h_i 's.
- A minibatch of N examples is randomly sampled and the contrastive prediction task is defined on pairs of augmented examples derived from the minibatch, resulting in 2N data points.
- Negative examples are not sampled explicitly. Instead, given a positive pair, the other 2(N-1) augmented examples within a minibatch are treated as negative examples.
- A NT-Xent (the normalized temperature-scaled cross entropy loss) loss function is used:

$$\ell_{i,j} = -\log rac{\exp(\mathrm{sim}(z_i,z_j)/ au)}{\sum_{k=1}^{2N} \mathbb{1}_{[k
eq i]} \exp(\mathrm{sim}(z_i,z_k)/ au)}$$

where
$$\mathrm{sim}(z_i,z_j) = rac{z_i^T z_j}{\|z_i\| \|z_j\|}$$

• PyTorch Code (https://github.com/sthalles/SimCLR)

Momentum Contrast (MoCo)

- **Momentum Contrast (MoCo)** (https://arxiv.org/abs/1911.05722) is a self-supervised learning algorithm with a contrastive loss.
- · Contrastive loss methods can be thought of as building dynamic dictionaries.
- The "keys" (tokens) in the dictionary are sampled from data (e.g., images or patches) and are represented by an encoder network.
- Unsupervised learning trains encoders (by minimizing a contrastive loss) to perform dictionary look-up: an encoded "query" should be similar to its matching key and dissimilar to others.
- In MoCo, we maintain the dictionary as a queue of data samples: the encoded representations of the current mini-batch are enqueued, and the oldest are dequeued.
- The queue decouples the dictionary size from the mini-batch size, allowing it to be large.
- Moreover, as the dictionary keys come from the preceding several mini-batches, a slowly progressing key encoder, implemented as a
 momentum-based moving average of the query encoder, is proposed to maintain consistency.
- PyTorch Code (https://github.com/facebookresearch/moco)
 - Colab Demo (https://colab.research.google.com/github/facebookresearch/moco/blob/colab-notebook/colab/moco_cifar10_demo.ipynb)

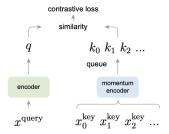


Figure 1. Momentum Contrast (MoCo) trains a visual representation encoder by matching an encoded query q to a dictionary of encoded keys using a contrastive loss. The dictionary keys $\{k_0,k_1,k_2,\ldots\}$ are defined on-the-fly by a set of data samples. The dictionary is built as a queue, with the current mini-batch enqueued and the oldest mini-batch dequeued, decoupling it from the mini-batch size. The keys are encoded by a slowly progressing encoder, driven by a momentum update with the query encoder. This method enables a large and consistent dictionary for learning visual representations.



Bootstrap Your Own Latent (BYOL)

- Bootstrap Your Own Latent (BYOL) (https://arxiv.org/abs/2006.07733) builds on the momentum network concept of MoCo, adding an MLP q_{θ} to predict z' from z.
- Rather than using a contrastive loss, BYOL uses the L2 error between the normalized prediction p and target z'.
- BYOL produces two augmented views $v \triangleq t(x)$ and $v' \triangleq t'(x)$ from x by applying respectively image augmentations $t \sim \mathcal{T}$ and $t' \sim \mathcal{T}'$.
- BYOL tries to convert both augmentaations of an image into the same representation vector (make p and z' equal).
- The L2 loss function does not require negative examples!
- There is an implicit contrastive loss, by using Batch Normalization in the first layers of the MLPs (MoCo does not require it).
- PyTorch Code (https://github.com/lucidrains/byol-pytorch)

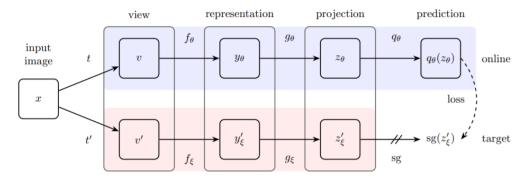
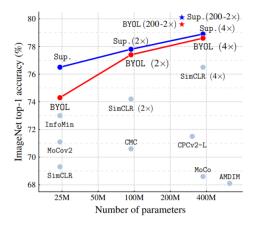


Figure 2: BYOL's architecture. BYOL minimizes a similarity loss between $q_{\theta}(z_{\theta})$ and $\operatorname{sg}(z'_{\xi})$, where θ are the trained weights, ξ are an exponential moving average of θ and sg means stop-gradient. At the end of training, everything but f_{θ} is discarded, and y_{θ} is used as the image representation.

- One purpose of negative examples in a contrastive loss function is to prevent *mode collapse*.
 - An example of mode collapse would be a network that always outputs $[1,0,0,0,\ldots]$ as its projection vector z.
 - If all projection vectors z are the same, then the network only needs to learn the identity function for q in order to achieve perfect prediction accuracy!
- If **batch normalization** is used in the projection layer g, the projection output vector z cannot collapse to any singular value like $[1,0,0,0,\ldots]$ because that is exactly what batch normalization prevents.
- Regardless of how similar the inputs to the batch normalization layer, the outputs will be redistributed according to the **learned** mean and standard deviation (and scaling-shifting).
- Mode collapse is prevented precisely because all samples in the mini-batch cannot take on the same value after batch normalization.
- · Said another way, with batch normalization, BYOL learns by asking, "how is this image different from the average image?".
 - The explicit contrastive approach used by SimCLR and MoCo learns by asking, "what distinguishes these two specific images from each other?"
 - These two approaches seem equivalent, since comparing an image with many other images has the same effect as comparing it to the average of the other images.

Performance Comparison

Performance on ImageNet (linear evaluation) using ResNet-50 and ResNet200 (2x), compared to other unsupervised and supervised (Sup.) baselines:







- These videos do not replace the lectures and tutorials.
- · Please use these to get a better understanding of the material, and not as an alternative to the written material.

Video By Subject

- Transfer Learning <u>Transfer Learning (C3W2L07) (https://www.youtube.com/watch?v=yofjFQddwHE)</u>
 - Transfer Learning in PyTorch-PyTorch Tutorial 15 Transfer Learning (https://www.youtube.com/watch?v=K0lWSB2QolQ)
- General Self-Supervised Learning <u>Lecture 7 Self-Supervised Learning -- UC Berkeley Spring 2020 CS294-158 Deep Unsupervised Learning (https://www.youtube.com/watch?v=dMUes74-nYY)</u>
- SimCLR SimCLR Explained! (https://www.youtube.com/watch?v=APki8LmdJwY)
- MoCo Momentum Contrastive Learning (https://www.youtube.com/watch?v=LvHwBQF14zs)
- BYOL BYOL: Bootstrap Your Own Latent: A New Approach to Self-Supervised Learning (https://www.youtube.com/watch?v=YPfUiOMYOEE)



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- Sebastian Ruder Transfer Learning Machine Learning's Next Frontier (https://ruder.io/transfer-learning/)
- Jacob Devlin and Ming-Wei Chang Open Sourcing BERT: State-of-the-Art Pre-training for Natural Language Processing (https://ai.googleblog.com/2018/11/open-sourcing-bert-state-of-art-pre.html)
- CS294-158-SP20-Deep Unsupervised Learning (https://sites.google.com/view/berkeley-cs294-158-sp20/home)
- Contrastive Predictive Coding (https://paperswithcode.com/method/contrastive-predictive-coding)
- Simple Framework for Contrastive Learning of Visual Representations (SimCLR) (https://paperswithcode.com/method/simclr)
- Momentum Contrast (https://paperswithcode.com/method/moco)
- Understanding self-supervised and contrastive learning with "Bootstrap Your Own Latent" (BYOL) (https://untitled-ai.github.io/understanding-self-supervised-contrastive-learning.html)