Style Transfer

In this notebook we will implement the style transfer technique from "Image Style Transfer Using Convolutional Neural Networks" (Gatys et al., CVPR 2015) (http://www.cv-

foundation.org/openaccess/content cvpr 2016/papers/Gatys Image Style Transfer CVPR 2016 paper.pdf).

The general idea is to take two images, and produce a new image that reflects the content of one but the artistic "style" of the other. We will do this by first formulating a loss function that matches the content and style of each respective image in the feature space of a deep network, and then performing gradient descent on the pixels of the image itself.

The deep network we use as a feature extractor is SqueezeNet (https://arxiv.org/abs/1602.07360), a small model that has been trained on ImageNet. You could use any network, but we chose SqueezeNet here for its small size and efficiency.

Here's an example of the images you'll be able to produce by the end of this notebook:

Style Source



Content Source



Output Image



Setup

```
In [1]: import torch
        import torch.nn as nn
        import torchvision
        import torchvision.transforms as T
```

import PIL

import numpy as np

from scipy.misc import imread from collections import namedtuple import matplotlib.pyplot as plt

from cs231n.image utils import SQUEEZENET MEAN, SQUEEZENET STD %matplotlib inline

We provide you with some helper functions to deal with images, since for this part of the assignment we're dealing with real JPEGs, not CIFAR-10 data.

```
In [2]: def preprocess(img, size=512):
            transform = T.Compose([
                T.Resize(size),
                T.ToTensor(),
                T.Normalize(mean=SQUEEZENET MEAN.tolist(),
                             std=SQUEEZENET_STD.tolist()),
                T.Lambda(lambda x: x[None]),
            ])
            return transform(img)
        def deprocess(img):
            transform = T.Compose([
                T.Lambda (lambda x: x[0]),
                T.Normalize(mean=[0, 0, 0], std=[1.0 / s for s in SQUEEZENET_STD
        .tolist()]),
                T.Normalize(mean=[-m for m in SQUEEZENET MEAN.tolist()], std=[1,
         1, 1]),
                T.Lambda(rescale),
                T.ToPILImage(),
            ])
            return transform(img)
        def rescale(x):
            low, high = x.min(), x.max()
            x rescaled = (x - low) / (high - low)
            return x rescaled
        def rel error(x,y):
            return np.max(np.abs(x - y) / (np.maximum(1e-8, np.abs(x) + np.abs(y)
        ))))
        def features from img(imgpath, imgsize):
            img = preprocess(PIL.Image.open(imgpath), size=imgsize)
            img var = img.type(dtype)
            return extract features(img var, cnn), img var
        # Older versions of scipy.misc.imresize yield different results
        # from newer versions, so we check to make sure scipy is up to date.
        def check scipy():
            import scipy
            vnum = int(scipy. version .split('.')[1])
            major vnum = int(scipy. version .split('.')[0])
            assert vnum >= 16 or major_vnum >= 1, "You must install SciPy >= 0.1
        6.0 to complete this notebook."
        check scipy()
        answers = dict(np.load('style-transfer-checks.npz'))
```

As in the last assignment, we need to set the dtype to select either the CPU or the GPU

In [3]: dtype = torch.FloatTensor
Uncomment out the following line if you're on a machine with a GPU set
 up for PyTorch!
#dtype = torch.cuda.FloatTensor

```
In [5]: # Load the pre-trained SqueezeNet model.
        cnn = torchvision.models.squeezenet1 1(pretrained=True).features
        cnn.type(dtype)
        # We don't want to train the model any further, so we don't want PyTorch
         to waste computation
        # computing gradients on parameters we're never going to update.
        for param in cnn.parameters():
            param.requires_grad = False
        # We provide this helper code which takes an image, a model (cnn), and r
        eturns a list of
        # feature maps, one per layer.
        def extract features(x, cnn):
            Use the CNN to extract features from the input image x.
            Inputs:
            - x: A PyTorch Tensor of shape (N, C, H, W) holding a minibatch of i
        mages that
              will be fed to the CNN.
            - cnn: A PyTorch model that we will use to extract features.
            Returns:
            - features: A list of feature for the input images x extracted using
         the cnn model.
              features[i] is a PyTorch Tensor of shape (N, C i, H i, W i); recal
        1 that features
              from different layers of the network may have different numbers of
         channels (C i) and
              spatial dimensions (H i, W i).
            features = []
            prev feat = x
            for i, module in enumerate(cnn._modules.values()):
                next feat = module(prev feat)
                features.append(next feat)
                prev feat = next feat
            return features
        #please disregard warnings about initialization
```

/home/shared/anaconda3/lib/python3.6/site-packages/torchvision-0.2.1-py 3.6.egg/torchvision/models/squeezenet.py:94: UserWarning: nn.init.kaiming_uniform is now deprecated in favor of nn.init.kaiming_uniform_. /home/shared/anaconda3/lib/python3.6/site-packages/torchvision-0.2.1-py 3.6.egg/torchvision/models/squeezenet.py:92: UserWarning: nn.init.norma l is now deprecated in favor of nn.init.normal_.

Computing Loss

We're going to compute the three components of our loss function now. The loss function is a weighted sum of three terms: content loss + style loss + total variation loss. You'll fill in the functions that compute these weighted terms below.

Content loss

We can generate an image that reflects the content of one image and the style of another by incorporating both in our loss function. We want to penalize deviations from the content of the content image and deviations from the style of the style image. We can then use this hybrid loss function to perform gradient descent **not on the parameters** of the model, but instead **on the pixel values** of our original image.

Let's first write the content loss function. Content loss measures how much the feature map of the generated image differs from the feature map of the source image. We only care about the content representation of one layer of the network (say, layer ℓ), that has feature maps $A^\ell \in \mathbb{R}^{1 \times C_\ell \times H_\ell \times W_\ell}$. C_ℓ is the number of filters/channels in layer ℓ , H_ℓ and W_ℓ are the height and width. We will work with reshaped versions of these feature maps that combine all spatial positions into one dimension. Let $F^\ell \in \mathbb{R}^{C_\ell \times M_\ell}$ be the feature map for the current image and $P^\ell \in \mathbb{R}^{C_\ell \times M_\ell}$ be the feature map for the content source image where $M_\ell = H_\ell \times W_\ell$ is the number of elements in each feature map. Each row of F^ℓ or P^ℓ represents the vectorized activations of a particular filter, convolved over all positions of the image. Finally, let w_ℓ be the weight of the content loss term in the loss function.

Then the content loss is given by:

$$L_c = w_c \times \sum_{i,j} (F_{ij}^{\ell} - P_{ij}^{\ell})^2$$

Test your content loss. You should see errors less than 0.0001.

Maximum error is 0.000

Style loss

Now we can tackle the style loss. For a given layer ℓ , the style loss is defined as follows:

First, compute the Gram matrix G which represents the correlations between the responses of each filter, where F is as above. The Gram matrix is an approximation to the covariance matrix -- we want the activation statistics of our generated image to match the activation statistics of our style image, and matching the (approximate) covariance is one way to do that. There are a variety of ways you could do this, but the Gram matrix is nice because it's easy to compute and in practice shows good results.

Given a feature map F^{ℓ} of shape (C_{ℓ}, M_{ℓ}) , the Gram matrix has shape (C_{ℓ}, C_{ℓ}) and its elements are given by:

$$G_{ij}^{\ell} = \sum_{k} F_{ik}^{\ell} F_{jk}^{\ell}$$

Assuming G^ℓ is the Gram matrix from the feature map of the current image, A^ℓ is the Gram Matrix from the feature map of the source style image, and w_ℓ a scalar weight term, then the style loss for the layer ℓ is simply the weighted Euclidean distance between the two Gram matrices:

$$L_s^{\ell} = w_{\ell} \sum_{i,j} \left(G_{ij}^{\ell} - A_{ij}^{\ell} \right)^2$$

In practice we usually compute the style loss at a set of layers \mathcal{L} rather than just a single layer ℓ ; then the total style loss is the sum of style losses at each layer:

$$L_s = \sum_{\ell \in \mathcal{C}} L_s^{\ell}$$

Begin by implementing the Gram matrix computation below:

```
In [9]: def gram_matrix(features, normalize=True):
            Compute the Gram matrix from features.
            Inputs:
            - features: PyTorch Tensor of shape (N, C, H, W) giving features for
              a batch of N images.
            - normalize: optional, whether to normalize the Gram matrix
                If True, divide the Gram matrix by the number of neurons (H * W
         * C)
            Returns:
            - gram: PyTorch Tensor of shape (N, C, C) giving the
              (optionally normalized) Gram matrices for the N input images.
            N, C, H, W = features.shape
            temp = features.view(N,C,-1)
            gram = torch.zeros((N,C,C))
            for i in range(N) :
                gram[i] = temp[i].matmul(temp[i].transpose(0,1))
            if (normalize) :
                gram /= H*W*C
            return gram
```

Test your Gram matrix code. You should see errors less than 0.0001.

```
In [11]: def gram_matrix_test(correct):
    style_image = 'styles/starry_night.jpg'
    style_size = 192
    feats, _ = features_from_img(style_image, style_size)
    student_output = gram_matrix(feats[5].clone()).cpu().data.numpy()
    error = rel_error(correct, student_output)
    print('Maximum error is {:.3f}'.format(error))
```

Maximum error is 0.000

Next, implement the style loss:

```
In [18]: # Now put it together in the style loss function...
         def style loss(feats, style layers, style targets, style weights):
             Computes the style loss at a set of layers.
             Inputs:
             - feats: list of the features at every layer of the current image, a
         s produced by
               the extract features function.
             - style layers: List of layer indices into feats giving the layers t
         o include in the
               style loss.
             - style targets: List of the same length as style layers, where styl
         e targets[i] is
               a PyTorch Tensor giving the Gram matrix of the source style image
          computed at
               layer style layers[i].
             - style weights: List of the same length as style layers, where styl
         e weights[i]
               is a scalar giving the weight for the style loss at layer style la
         yers[i].
             Returns:
             - style loss: A PyTorch Tensor holding a scalar giving the style los
         s.
             # Hint: you can do this with one for loop over the style layers, and
          should
             # not be very much code (~5 lines). You will need to use your gram m
         atrix function.
             loss = 0
             for i in range(len(style_layers)):
                 loss += ((style targets[i] - gram matrix(feats[style layers[i
         ]]))**2).sum() * style weights[i]
             return loss
```

Test your style loss implementation. The error should be less than 0.0001.

```
In [19]: def style_loss_test(correct):
             content_image = 'styles/tubingen.jpg'
             style_image = 'styles/starry_night.jpg'
             image_size = 192
             style_size = 192
             style_layers = [1, 4, 6, 7]
             style_weights = [300000, 1000, 15, 3]
             c_feats, _ = features_from_img(content_image, image size)
             feats, _ = features_from_img(style_image, style_size)
             style_targets = []
             for idx in style_layers:
                 style_targets.append(gram_matrix(feats[idx].clone()))
             student_output = style_loss(c_feats, style_layers, style_targets, st
         yle_weights).cpu().data.numpy()
             error = rel_error(correct, student_output)
             print('Error is {:.3f}'.format(error))
         style loss test(answers['sl out'])
```

Error is 0.000

Total-variation regularization

It turns out that it's helpful to also encourage smoothness in the image. We can do this by adding another term to our loss that penalizes wiggles or "total variation" in the pixel values.

You can compute the "total variation" as the sum of the squares of differences in the pixel values for all pairs of pixels that are next to each other (horizontally or vertically). Here we sum the total-variation regualarization for each of the 3 input channels (RGB), and weight the total summed loss by the total variation weight, W_t :

$$L_{tv} = w_t \times \left(\sum_{c=1}^3 \sum_{i=1}^{H-1} \sum_{j=1}^W (x_{i+1,j,c} - x_{i,j,c})^2 + \sum_{c=1}^3 \sum_{i=1}^H \sum_{j=1}^{W-1} (x_{i,j+1,c} - x_{i,j,c})^2\right)$$

In the next cell, fill in the definition for the TV loss term. To receive full credit, your implementation should not have any loops.

Test your TV loss implementation. Error should be less than 0.0001.

```
In [22]: def tv_loss_test(correct):
    content_image = 'styles/tubingen.jpg'
    image_size = 192
    tv_weight = 2e-2

    content_img = preprocess(PIL.Image.open(content_image), size=image_s
    ize)

    student_output = tv_loss(content_img, tv_weight).cpu().data.numpy()
    error = rel_error(correct, student_output)
    print('Error is {:.3f}'.format(error))

tv_loss_test(answers['tv_out'])
```

Error is 0.000

Now we're ready to string it all together (you shouldn't have to modify this function):

```
def style transfer(content image, style image, image size, style size, c
ontent layer, content weight,
                   style_layers, style_weights, tv_weight, init_random =
 False):
    11 11 11
    Run style transfer!
    Inputs:
    - content image: filename of content image
    - style image: filename of style image
    - image size: size of smallest image dimension (used for content los
s and generated image)
    - style size: size of smallest style image dimension
    - content layer: layer to use for content loss
    - content weight: weighting on content loss
    - style layers: list of layers to use for style loss
    - style weights: list of weights to use for each layer in style laye
rs
    - tv weight: weight of total variation regularization term
    - init random: initialize the starting image to uniform random noise
    # Extract features for the content image
    content img = preprocess(PIL.Image.open(content image), size=image s
ize)
    feats = extract_features(content_img, cnn)
    content target = feats[content layer].clone()
    # Extract features for the style image
    style img = preprocess(PIL.Image.open(style image), size=style size)
    feats = extract features(style img, cnn)
    style targets = []
    for idx in style layers:
        style targets.append(gram matrix(feats[idx].clone()))
    # Initialize output image to content image or nois
    if init random:
        img = torch.Tensor(content img.size()).uniform (0, 1).type(dtype
)
    else:
        img = content_img.clone().type(dtype)
    # We do want the gradient computed on our image!
    img.requires grad ()
    # Set up optimization hyperparameters
    initial lr = 3.0
    decayed_lr = 0.1
    decay lr at = 180
    # Note that we are optimizing the pixel values of the image by passi
ng
    # in the img Torch tensor, whose requires grad flag is set to True
    optimizer = torch.optim.Adam([img], lr=initial lr)
    f, axarr = plt.subplots(1,2)
```

```
axarr[0].axis('off')
    axarr[1].axis('off')
    axarr[0].set_title('Content Source Img.')
    axarr[1].set title('Style Source Img.')
    axarr[0].imshow(deprocess(content img.cpu()))
    axarr[1].imshow(deprocess(style_img.cpu()))
    plt.show()
    plt.figure()
    for t in range(200):
        if t < 190:
            img.data.clamp_(-1.5, 1.5)
        optimizer.zero_grad()
        feats = extract_features(img, cnn)
        # Compute loss
        c_loss = content_loss(content_weight, feats[content_layer], cont
ent_target)
        s loss = style loss(feats, style layers, style targets, style we
ights)
        t_loss = tv_loss(img, tv_weight)
        loss = c_loss + s_loss + t_loss
        loss.backward()
        # Perform gradient descents on our image values
        if t == decay lr at:
            optimizer = torch.optim.Adam([img], lr=decayed lr)
        optimizer.step()
        if t % 100 == 0:
            print('Iteration {}'.format(t))
            plt.axis('off')
            plt.imshow(deprocess(img.data.cpu()))
            plt.show()
    print('Iteration {}'.format(t))
    plt.axis('off')
    plt.imshow(deprocess(img.data.cpu()))
    plt.show()
```

Generate some pretty pictures!

Try out style_transfer on the three different parameter sets below. Make sure to run all three cells. Feel free to add your own, but make sure to include the results of style transfer on the third parameter set (starry night) in your submitted notebook.

- The content image is the filename of content image.
- The style image is the filename of style image.
- The image_size is the size of smallest image dimension of the content image (used for content loss and generated image).
- The style_size is the size of smallest style image dimension.
- The content_layer specifies which layer to use for content loss.
- The content_weight gives weighting on content loss in the overall loss function. Increasing the value of this parameter will make the final image look more realistic (closer to the original content).
- style layers specifies a list of which layers to use for style loss.
- style_weights specifies a list of weights to use for each layer in style_layers (each of which will
 contribute a term to the overall style loss). We generally use higher weights for the earlier style layers
 because they describe more local/smaller scale features, which are more important to texture than
 features over larger receptive fields. In general, increasing these weights will make the resulting image
 look less like the original content and more distorted towards the appearance of the style image.
- tv_weight specifies the weighting of total variation regularization in the overall loss function.
 Increasing this value makes the resulting image look smoother and less jagged, at the cost of lower fidelity to style and content.

Below the next three cells of code (in which you shouldn't change the hyperparameters), feel free to copy and paste the parameters to play around them and see how the resulting image changes.

```
In [25]: # Composition VII + Tubingen
    params1 = {
        'content_image' : 'styles/tubingen.jpg',
        'style_image' : 'styles/composition_vii.jpg',
        'image_size' : 192,
        'style_size' : 512,
        'content_layer' : 3,
        'content_weight' : 5e-2,
        'style_layers' : (1, 4, 6, 7),
        'style_weights' : (20000, 500, 12, 1),
        'tv_weight' : 5e-2
}
style_transfer(**params1)
```

Content Source Img.





Iteration 0



Iteration 100



Iteration 199



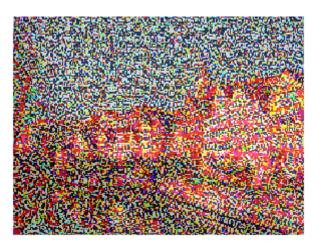
```
In [26]: # Scream + Tubingen
    params2 = {
        'content_image':'styles/tubingen.jpg',
        'style_image':'styles/the_scream.jpg',
        'image_size':192,
        'style_size':224,
        'content_layer':3,
        'content_weight':3e-2,
        'style_layers':[1, 4, 6, 7],
        'style_weights':[2000000, 800, 12, 1],
        'tv_weight':2e-2
}
style_transfer(**params2)
```

Content Source Img.





Iteration 0



Iteration 100



Iteration 199

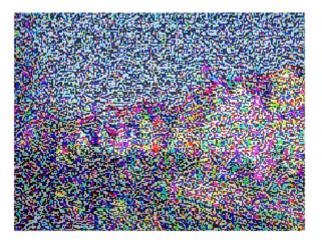


Content Source Img.





Iteration 0



Iteration 100



Iteration 199



Feature Inversion

The code you've written can do another cool thing. In an attempt to understand the types of features that convolutional networks learn to recognize, a recent paper [1] attempts to reconstruct an image from its feature representation. We can easily implement this idea using image gradients from the pretrained network, which is exactly what we did above (but with two different feature representations).

Now, if you set the style weights to all be 0 and initialize the starting image to random noise instead of the content source image, you'll reconstruct an image from the feature representation of the content source image. You're starting with total noise, but you should end up with something that looks quite a bit like your original image.

(Similarly, you could do "texture synthesis" from scratch if you set the content weight to 0 and initialize the starting image to random noise, but we won't ask you to do that here.)

Run the following cell to try out feature inversion.

[1] Aravindh Mahendran, Andrea Vedaldi, "Understanding Deep Image Representations by Inverting them", CVPR 2015

```
In [28]: # Feature Inversion -- Starry Night + Tubingen
    params_inv = {
        'content_image' : 'styles/tubingen.jpg',
        'style_image' : 'styles/starry_night.jpg',
        'image_size' : 192,
        'style_size' : 192,
        'content_layer' : 3,
        'content_weight' : 6e-2,
        'style_layers' : [1, 4, 6, 7],
        'style_weights' : [0, 0, 0, 0], # we discard any contributions from style to the loss
        'tv_weight' : 2e-2,
        'init_random': True # we want to initialize our image to be random
    }
    style_transfer(**params_inv)
```

Content Source Img.





Iteration 0



Iteration 100



Iteration 199

