

Style Transfer by Deep Learning

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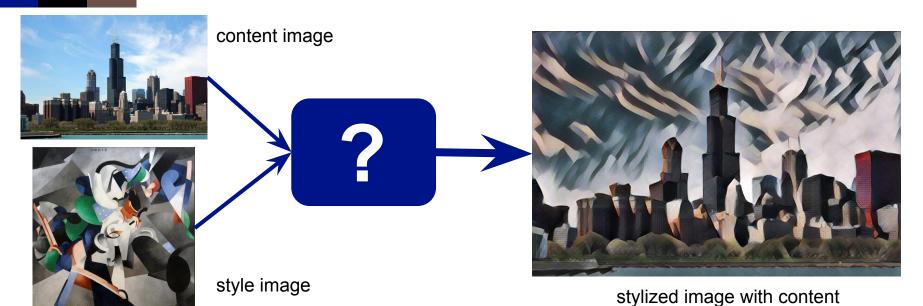








Style Transfer by Deep Learning: Motivations



- Generative task
 - From an image, generate a new one
- Introduction to more complex tasks
 - Super-resolution and colorisation
- CNNs understanding is required
 - Hierarchy of representations
 - Feature spaces



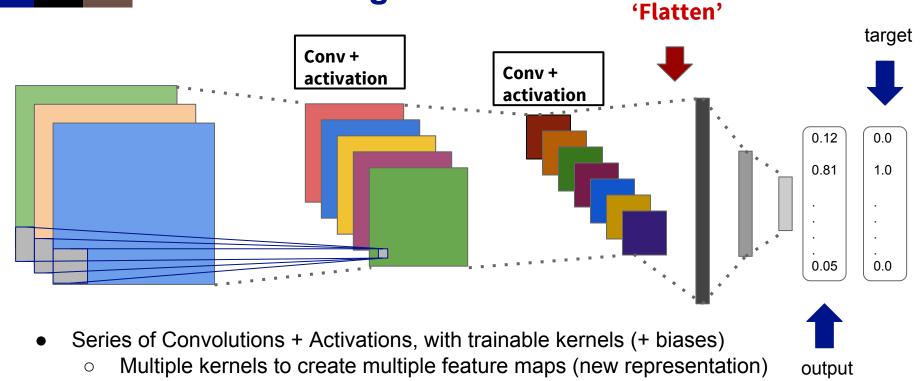
- CNNs visualizing and understanding
- Content and style representations
- Optimization-based style transfer
- Feed-forward method with learning
- Arbitrary style transfer



CNNs visualizing and understanding



CNNs for Image Classification

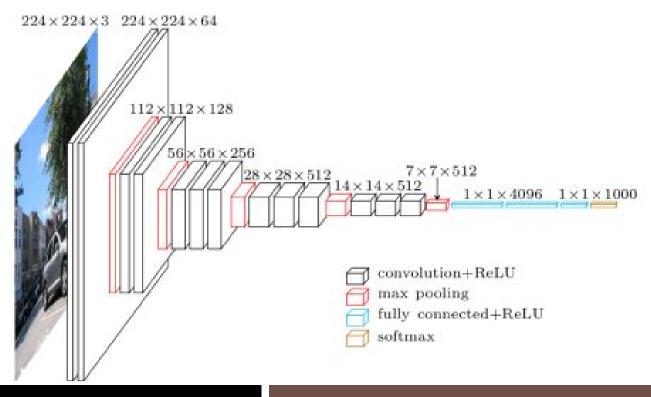


- Pooling operations to reduce the dimensions of the feature maps
- 'Flatten' operation, representation as a vector
- Fully-Connected layers (Multilayer Perceptron)
 - Learned: weight matrix and bias vector
- Training (weight-update) on error
 - Classification : cross entropy

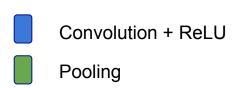


Deep Network: VGG-16

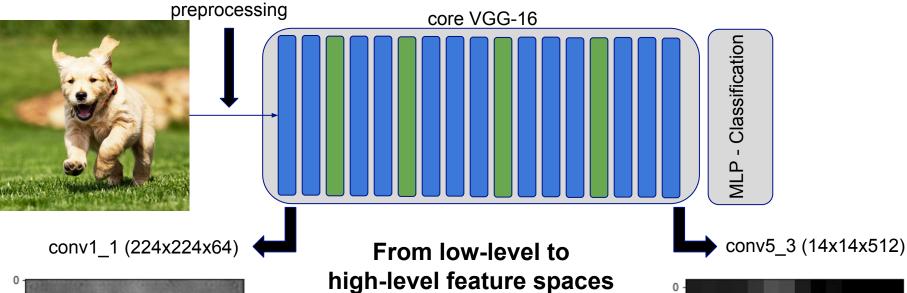
- Simple : no inception modules or residual connections
- Trained for **image classification** on ImageNet (1000 classes)
- State of the art in 2014 (92.7% top-5 test accuracy)
- 138,357,544 parameters (10% conv weights, 90% FC layers)
- No residual connections, or inception modules : Deep simple model

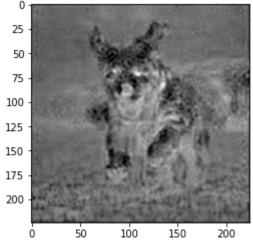






CNNs visualizing







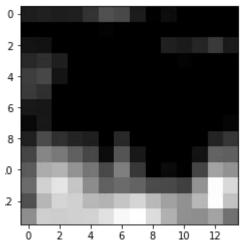






Additional visualization methods:

- Deep Dream approach [2]
 - Optimization-based
- Zeiler & Fergus [3]
 - Transposed convolutions and unpooling operations



Content & style representations



Content Representation/Reconstruction

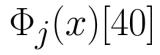
 x_c

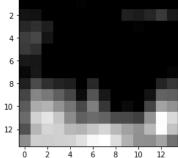


 $\Phi_j(x)$

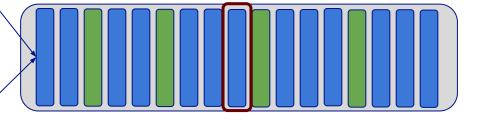
Activations of the jth layer

Eg:





conv3_3:56x56x256



Fixed VGG-16

$$\mathcal{X}$$

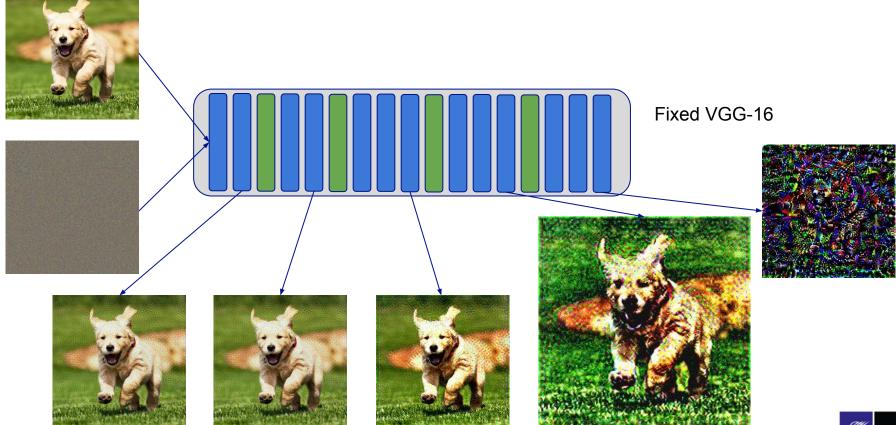
$$\hat{x} = \underset{x}{\operatorname{argmin}} ||\Phi_j(x) - \Phi_j(x_c)||_2^2$$

- Goal: find an image with the same activations at a given layer (all feature maps)
- Optimization problem, start from a random image



Content Representation/Reconstruction

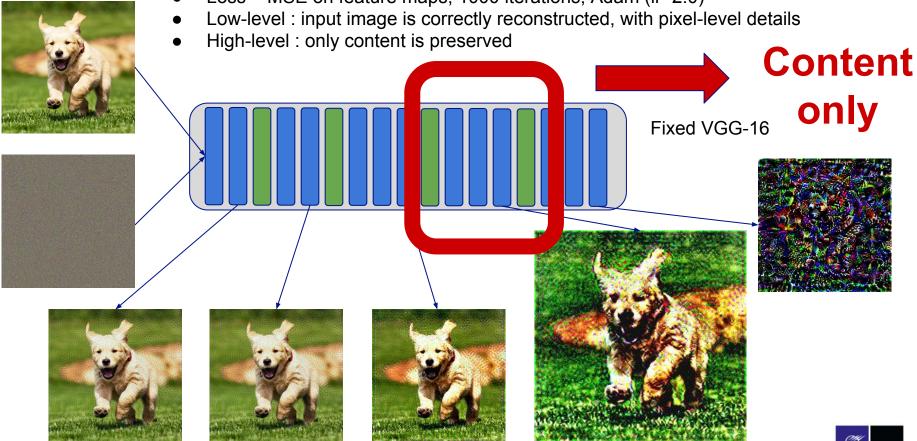
- gradient descent optimization on input image, network does not change
- loss = MSE on feature maps, 1000 iterations, Adam (Ir=2.0)
- low-level : input image is correctly reconstructed, with pixel-level details
- high-level : only content is preserved





Content Representation/Reconstruction

- From a random image, reconstruct the feature maps obtained with a normal image, on a specific layer
- Gradient descent optimization on image input, network does not change
 - Loss = MSE on feature maps, 1000 iterations, Adam (Ir=2.0)





Style Representation/Reconstruction

- Needs more complex statistics on feature maps : **Gram matrix**
 - Second-order statistics
 - Can capture texture information, no spatial information
- $\bullet \quad$ For a given layer j with C_j feature maps of size (W_j, H_j)
- ullet The Gram matrix is a (C_j,C_j) matrix :

$$G_j(x)_{c_1,c_2} = \mathbb{E}[\Phi_j(x)[c_1] * \Phi_j(x)[c_2]]$$

- Where * is an element-wise operation between 2 feature maps (Hadamard product)
- Contains the correlation between every pair of feature maps



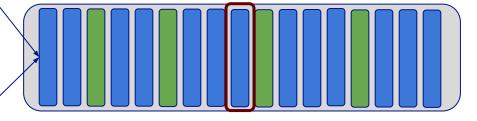
Style Representation/Reconstruction

 x_s

 $G_j(x)$

Gram matrix of the jth layer (256 x 256)

conv3_3:56x56x256



Fixed VGG-16

$$\overline{x}$$

$$\hat{x} = \underset{x}{\operatorname{argmin}} ||G_j(x) - G_j(x_s)||_2^2$$

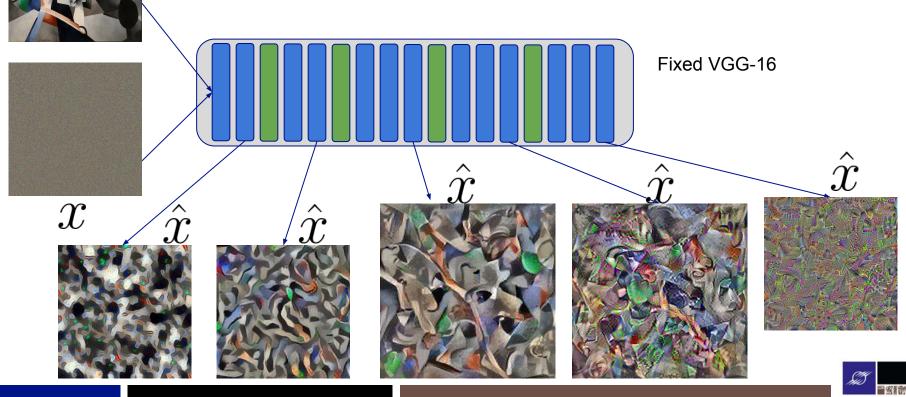
- Goal: To find an image with the same Gram matrix for a given layer
- Optimization problem: Start from a random image



Style Representation/Reconstruction

 \mathcal{X}_{S}

- Gradient descent optimization on image input, network is freezed
- Loss = MSE on feature maps, 1000 iterations, Adam (Ir=2.0)
- Low-level : Small and simple patterns
- High-level : More complex patterns



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Content & Style Representations

- Content is preserved in high level features
- Style is present in second-order statistics in low and medium levels
- Content and Style are separable

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- content_loss and a style_loss are defined
- Combine style and loss from different images is possible, via feature extraction learned within a VGG network, trained for image classification



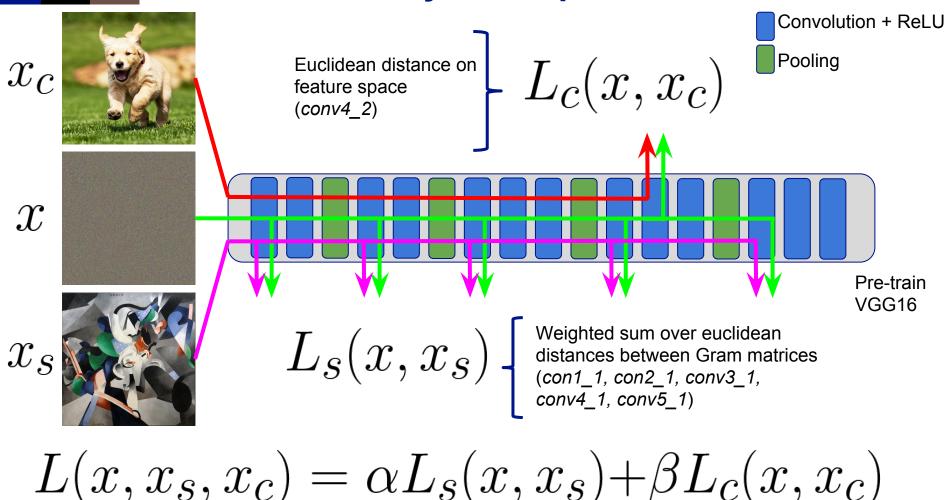
Optimization-based approach for style transfer

Approach proposed by

- Gatys et al [4, 10]
- Ruder et al [5]



Mix content & style via specific losses





Optimization process

- Compute content_target (feature maps) with content_image
- Compute style_target (Gram matrices) with style_image
- Start from a random image (input_image)
- Optimization process :
 - Compute content_loss and style_loss with targets + input_image
 - Minimize this loss by modifying input_image
 - Possible thanks to gradient-descent method (like Adam)



- TensorFlow implementation (version 1.1.0, Python 3.5)
- With TensorBoard annotations (Graph and metrics visualization)
- Jupyter notebooks and Conda/Docker envs
- GitHub: JGuillaumin/style-transfer-workshop



```
Content loss
with tf.name scope('content image'):
    # Construct content loss using content image.
    sess.run(vgg['input'].assign(content image))
with tf.name scope('content loss'):
    content target = sess.run(vgg['conv4 2'])
    N = content target.shape[3] # number of feature maps
   M = content target.shape[1] * content target.shape[2] # number of feature per feature map
    content loss = (1 / (4 * N * M)) * tf.reduce sum(tf.pow(vgg['conv4 2'] - content target, 2))
      Style loss
    with tf.name scope('style image'):
         # Construct style loss using style image
         sess.run(vgg['input'].assign(style image))
    STYLE_LAYERS = [('conv1_1', 0.5), ('conv2_1', 1.0), ('conv3_1', 2.5),
```

('conv4 1', 3.0), ('conv5 1', 1.0)]



```
def gram matrix tf(F, N, M):
    F = tf.reshape(F, (M, N))
    return tf.matmul(tf.transpose(F), F)
def gram matrix np(F, N, M):
    F = np.reshape(F, (M, N))
    return np.matmul(np.transpose(F), F)
with tf.name scope('style loss'):
    style loss = 0
    for layer name, weight in STYLE LAYERS:
        style target = sess.run(vgg[layer name])
        N = style target.shape[3] # number of feature maps
        M = style target.shape[1] * style target.shape[2] # number of features per feature map
        # compute Gram matrices : target and tensor
        style target = gram matrix np(style target, N, M) # works on Numpy array
        G = gram matrix tf(vgg[layer name], N, M) # works on Tensor
        style loss += weight * (1 / (4 * N**2 * M**2)) * tf.reduce sum(tf.pow(G - style target, 2))
```



```
with tf.name_scope('total_loss'):
    total_loss = BETA * content_loss + ALPHA * style_loss

with tf.name_scope('train'):
    optimizer = tf.train.AdamOptimizer(2.0)
    train_step = optimizer.minimize(total_loss)

_ = sess.run(vgg['input'].assign(noise_image))
```



```
ITERATIONS = 1000

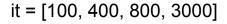
for it in range(ITERATIONS):
    _ = sess.run(train_step)

if it%100 == 0:
    _image = sess.run(vgg['input'])
    filename = 'output/stylized_gatys_iter{}.png'.format(it)
    save_image(filename, _image)
```











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Results

- Produce high-quality images
- Easy to tune effects (more content? more style?)
- Any input/output size
- Running time (1000 #iter)
 - GPU (GTX 1070): ~ 5 min (1920 CUDA cores)
 - CPU (i7-7700K): ~ 150 min (4 cores x 2 threads)
- Avoid any real-time applications
- But perceptual loss (content+style) is defined



Improvements

- Time dependency for video transformation (see [5])
- Change optimizer: L-BFGS!
- Tune weights between style and content loss
- Start from : content image, style image, noisy image, or a mix.
- Color constraint : preserve color from content image ! (see [10])



from : github.com/tensorflow/magenta



Feed-forward method with learning

Approach proposed by

- Ulyanov et al [6, 7]
- Johnson et al [8]



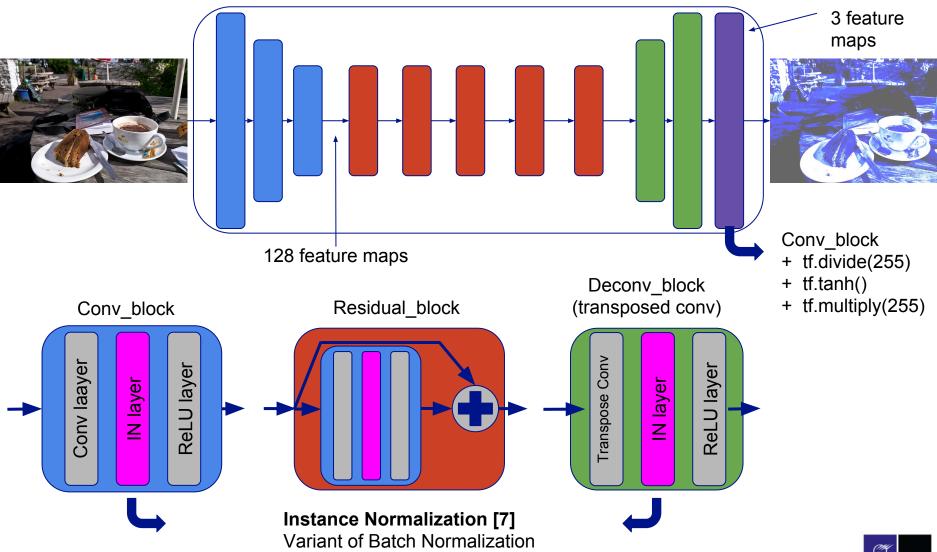
Feed-forward method

- Train a network to obtain a stylized image in one pass as an output
- Used for one specific style (fixed)

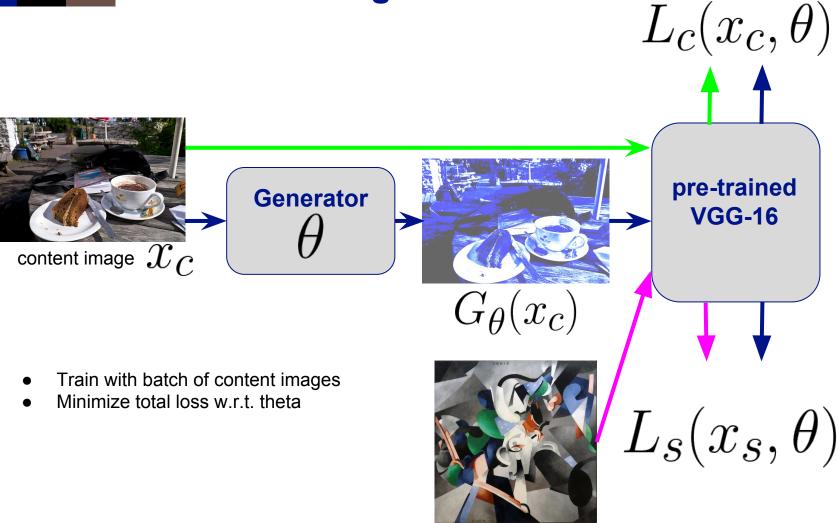




What type of structures for the generator?



How to train a generator?





Need a dataset of content images

- COCO dataset, about 80k images
- Only 1 style image

Training process (loop):

- Take a batch of samples from COCO
- Pass this batch through the generator to get generated images
- Compute style_loss between the generated images and the style image
- Compute content_loss between the generated images and the original ones
- Minimize the total_loss by updating the weights from the generator

Training information:

- Adam optimizer (*Ir=0.05*)
- Only 20k iterations (with batch_size=4)
- For 512x512x3:
 - Training time (on GTX 1070): 10 hours
 - Inference time: 330 ms (GTX 1070)



```
sess = tf.InteractiveSession()
generator={}
with tf.name scope('input images'):
    input images = tf.placeholder(tf.float32, shape=(BATCH SIZE, IMAGE HEIGHT, IMAGE WIDTH, 3),
                                 name='input images')
with tf.name scope('generator'):
    generator['conv block1'] = conv block(input images, 9, 32, 1, 'conv block1')
    generator['conv block2'] = conv block(generator['conv block1'], 3, 64, 2, 'conv block2')
    generator['conv block3'] = conv block(generator['conv block2'], 3, 128, 2, 'conv block3')
   generator['residual block1'] = residual block(generator['conv block3'], 'residual block1')
    generator['residual block2'] = residual block(generator['residual block1'], 'residual block2')
    generator['residual block3'] = residual block(generator['residual block2'], 'residual block3')
    generator['residual block4'] = residual block(generator['residual block3'], 'residual block4')
    generator['residual block5'] = residual block(generator['residual block4'], 'residual block5')
    generator['deconv block1'] = deconv block(generator['residual block5'], 3, 64, 2, 'deconv block1')
    generator['deconv block2'] = deconv block(generator['deconv block1'], 3, 32, 2, 'deconv block2')
    generator['final conv'] = conv block(generator['deconv block2'], 9, 3, 1, 'final conv', relu=False)
   generator['output'] = tf.multiply(tf.tanh(generator['final conv']/255.0), 255, name="output")
```



```
with tf.variables scope('VGGs', reuse=True):
    vgg = VGG.generate model(weights file=MODEL WEIGHTS,
                             input=generator['output'],
                             remove top=True,
                             with preprocessing=False)
    vgg content = VGG.generate model(weights file=MODEL WEIGHTS,
                                     input=input images,
                                      remove top=True,
                                     with preprocessing=False)
# variable sharing between 'vgg' and 'vgg content'
# Content loss
with tf.name scope('content loss'):
    content loss = tf.reduce mean(tf.pow(vgg content['conv4 2'] - vgg['conv4 2'], 2))
```





```
with tf.name scope('style loss'):
    style loss = 0
    for layer name, weight in STYLE LAYERS :
        style target = sess.run(vgg[layer name], feed dict=feed dict)
        B = style target.shape[0] # batch size
        N = style target.shape[3] # number of feature maps
        M = style target.shape[1] * style target.shape[2] # number of features per feature map
        # compute Gram matrices : target and tensor
        style target = gram matrix np(style target, B, N, M) # works on Numpy array
        G = gram matrix tf(vgg[layer name],B , N, M) # works on Tensor
        style loss += weight * tf.reduce mean(tf.pow(G - style target, 2))
with tf.name scope('total loss'):
    total loss = BETA * content loss + ALPHA * style loss
```



```
sess.run(tf.global_variables_initializer())
ITERATIONS = 20000
feed={}

for it in range(ITERATIONS):
    batch = COCO_batch_generator.next()
    feed[input_images] = batch

    _ = sess.run([train_step], feed_dict=feed)

if it%500 == 0:
    _image = sess.run(generator['output'], feed_dict={input_images:[content_image[0]]*BATCH_SIZE})
    save_image(filename, _image)
```



Results and improvements

it = 20000

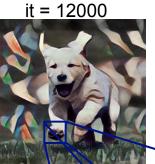
With a new content image :

it = 500

it = 1









- Learn to apply only one style!
- In [9] (ICLR 2017):
 - Add 'Conditional Instance Normalization'
 - Learn to apply a fixed set of styles (until 64)
 - Can learn quickly a new style (incremental learning)
- Use resized convolutions instead of transposed convolutions: Improves quality
- Add variational loss to encourage spatial smoothness



Arbitrary style transfer

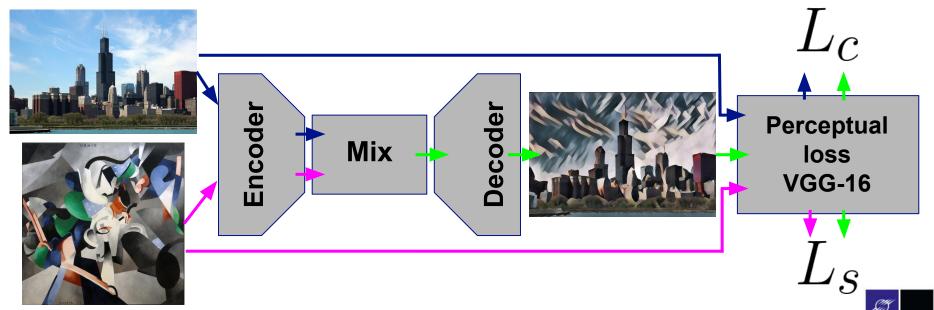
Approach proposed by:

- X. Huang, S. Belongie [11] (ICLR 2017)



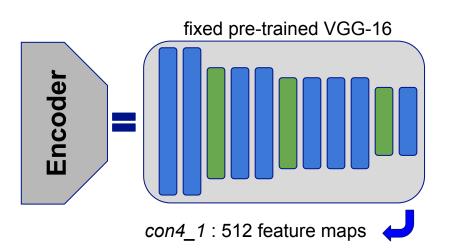
Mix content & style images within the generator

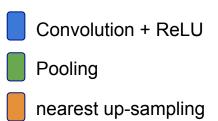
- Goal: train a generator to produce stylized images from any style with any content
- **Previous approach**: style is learned in the generator via the *style_loss*
- New approach : mix content and style images in feature space !
 - Use encoder-decoder structure
 - Mix encoded content and style images
 - Use the same perceptual loss (content_loss + style_loss)
 - Trained on a content dataset (COCO) and a painting dataset (WikiArt)

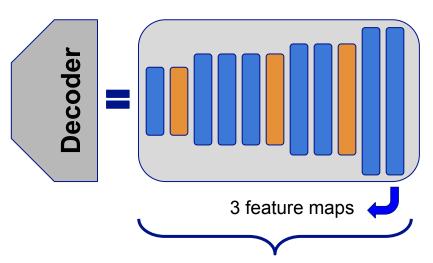


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Encoder/Decoder + Adaptive IN



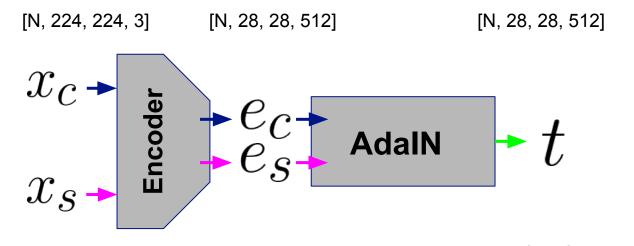




- the decoder mirrors the encoder
- to train!



Adaptive Instance Normalization

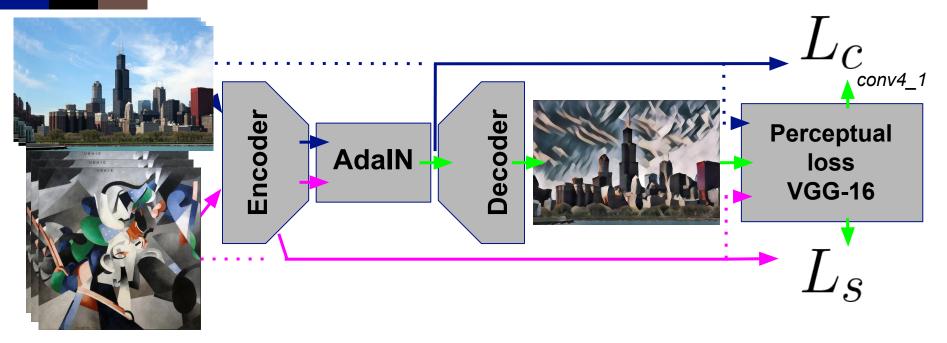


$$AdaIN_{n,f}(e_c, e_s) = \sigma_{n,f}(e_s) \frac{e_c - \mu_{n,f}(e_c)}{\sigma_{n,f}(e_c)} + \mu_{n,f}(e_s)$$

- per sample (n) and per channel (f) statistics alignment
- producing the target feature maps



How to train the decoder?



- trained with batches of content-styles image pairs
- different contents and styles within the same batch!
- 80k content images (MS COCO) + 80k paintings (WikiArt.org)



```
sess = tf.InteractiveSession()
generator={}
with tf.name scope('input images'):
    input content = tf.placeholder(tf.float32, shape=(BATCH SIZE, IMAGE HEIGHT, IMAGE WIDTH, 3),
                                  name='input content')
    input style= tf.placeholder(tf.float32, shape=(BATCH SIZE, IMAGE HEIGHT, IMAGE WIDTH, 3),
                                  name='input style')
with tf.variables scope('VGGs', reuse=True):
    content encoder = VGG.generate model(weights file=MODEL WEIGHTS,
                             input=input content,
                             remove top=True,
                             with preprocessing=False)
    style encoder = VGG.generate model(weights file=MODEL WEIGHTS,
                                     input=input style,
                                      remove top=True,
                                     with preprocessing=False)
encoded content = content encoder['conv4 1']
encoded style = style encoder['conv4 1']
```

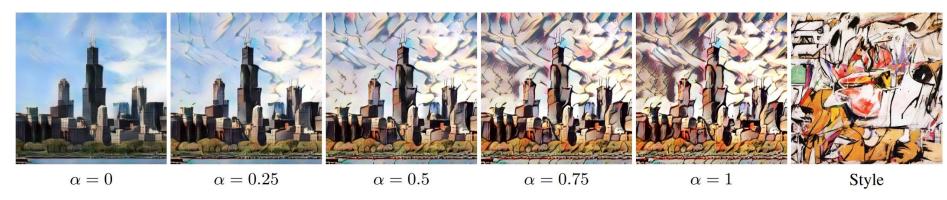
```
with tf.name scope('AdaIn') :
    eps = 1e-6
   mean c, var c = tf.nn.moments(encoded content, [1, 2], keep dims=True)
   mean s, var s = tf.nn.moments(encoded style, [1, 2], keep dims=True)
    target = mean s * ((encoded content - mean c)/(tf.sqrt(var c) + eps)) + mean s
with tf.name scope('decoder'):
    decoder = {}
    decoder['conv block1 1'] = conv block(taget, 3, 256 , 1, 'conv block1 1')
    decoder['up sampling1'] = up sampling(decoder['conv block1 1'], 2, 'up sampling1')
    decoder['conv block2 1'] = conv block(decoder['up sampling1'], 3, 256 , 1, 'conv block2 1')
    decoder['conv block2 2'] = conv block(decoder['conv block2 1'], 3, 256 , 1, 'conv block2 2')
    decoder['conv block2 3'] = conv block(decoder['conv block2 2'], 3, 128 , 1, 'conv block2 3')
    decoder['up sampling2'] = up sampling(decoder['conv block2 3'], 2, 'up sampling2')
    decoder['conv block3 1'] = conv block(decoder['up sampling2'], 3, 128 , 1, 'conv block3 1')
    decoder['conv block3 2'] = conv block(decoder['conv block3 1'], 3, 64 , 1, 'conv block3 2')
    decoder['up sampling3'] = up sampling(decoder['conv block3 2'], 2, 'up sampling3')
    decoder['conv block4 1'] = conv block(decoder['up sampling3'], 3, 64 , 1, 'conv block3 1')
    decoder['conv block4 2'] = conv block(decoder['conv block4 1'], 3, 64 , 1, 'conv block4 2')
    decoder['final conv'] = conv block(decoder['conv block4 1'], 9, 3, 1, 'final conv', relu=False)
    decoder['output'] = tf.multiply(tf.tanh(decoder['final conv']/255.0), 255, name="output")
```

```
with tf.variables scope('VGGs', reuse=True):
    vgg = VGG.generate model(weights file=MODEL WEIGHTS,
                             input=decoder['output'],
                             remove top=True,
                             with preprocessing=False)
with tf.name scope('content loss'):
    content loss = tf.reduce mean(tf.pow(vgg['conv4 1'] - target, 2))
with tf.name scope('style loss'):
    style loss = 0
    for layer name, weight in STYLE LAYERS :
        shape = vgg[layer name].get shape().as list()
        B = style target.shape[0] # batch size
        N = style target.shape[3] # number of feature maps
        M = style target.shape[1] * style target.shape[2] # number of features per feature map
        G_style = _gram_matrix tf(style encoder[layer name], B, N, M) # works on Numpy array
        G = gram matrix tf(vgg[layer name], B , N, M) # works on Tensor
        style loss += weight * tf.reduce mean(tf.pow(G - G style, 2))
```

```
with tf.name scope('total loss'):
    total loss = BETA * content loss + ALPHA * style loss
with tf.name scope('train'):
    optimizer = tf.train.AdamOptimizer(0.02)
    train step = optimizer.minimize(total loss)
sess.run(tf.global variables initializer())
ITERATIONS = 20000
feed={}
for it in range(ITERATIONS):
    batch c = COCO batch generator.next()
    batch s = WikiArt batch generator.next()
    feed[input content] = batch c
    feed[input style] = batch s
      = sess.run([train step], feed dict=feed)
    if it%500 == 0:
        image = sess.run(generator['output'], feed dict={input content:[content image[0]]*BATCH SIZE,
                                                          input style:[style image[0]]*BATCH SIZE})
        save image(filename, image)
```

Results

- Training time: about 24 hours (GTX 1070, 200k iterations)
- Decoder-AdalN-decoder can apply any new style!
- High-quality images (similar to optimization-based approach)
- Inference time : ~ 450ms (256x256)



from github.com/xunhuang1995/AdalN-style



Conclusion - Style Transfer

	Content	Style	Production	Training
Gatys et al (2015)	1	1	~ 5 min	X
Johnson Ulyanov (2016)	infinite	fixed number, until 64	~ 300ms	~ 4 hours
X Huang (2017)	infinite	infinite	~ 450 ms	~ days



Conclusion

- Resolve complex task by working on feature spaces
- Introduction to more complex tasks
 - Colorisation
 - Super-resolution
 - Inverse style transfer
 - Season transfer
 - Color transfer





from github.com/junyanz/CycleGAN



Online resources



github.com/JGuillaumin/style-transfer-workshop

- Several Jupyter notebooks
- All methods presented here and more
- Implementation with TensorFlow 1.1, Python 3.5
- With TensorBoard annotations
- Conda env and Dockerfiles (CPU and GPU)

Online in few days



Thank you



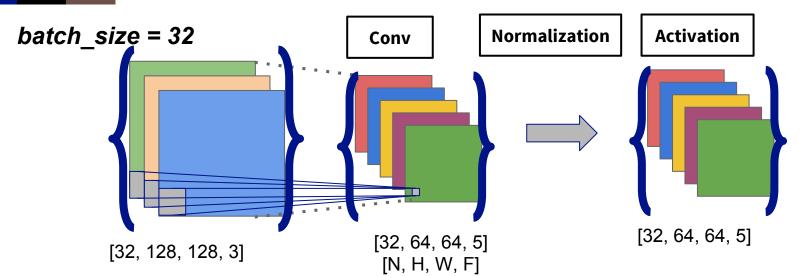


Resources

- [1]: K. Simonyan, A. Zisserman: "Very Deep Convolutional Networks for Large-Scale Image Recognition", 2014, arXiv:1409.1556
- [2] : About Deep Dream visualization technique : "Inceptionism: Going Deeper into Neural Networks"
- [3] : M. Zeiler, R. Fergus: Visualizing and Understanding Convolutional Networks, 2013 arXiv:1311.2901
- [4]: L. Gatys, A. Ecker, M. Bethge: A neural algorithm of artistic style, 2015, arXiv:1508.06576
- [5]: M. Ruder, A. Dosovitskiy, T. Brox: Artistic style transfer for video, 2016, arXiv:1604.08610
- [6]: D. Ulyanov et al: Texture Networks: Feed-forward Synthesis of Textures and Stylized Images, 2016, arXiv:1603.03417
- [7]: D. Ulyanov et al: Instance Normalization: The Missing Ingredient for Fast Stylization, 2016, arXiv:1607.08022
- [8] : J. Johnson et al : Perceptual losses for real-time style transfer and super-resolution, 2016, arXiv:1603.08155
- [9]: V. Dumoulin et al: A learned representation for artistic style, 2017, arXiv:1610.07629
- [10] : Gatys et al : Preserving color in Neural Artistic Style Transfer, 2016, arXiv:1606.05897
- [11] X. Huang and S. Belongie: Arbitrary Style Transfer in real-time with AdalN, 2017, arXiv:1703.06868



Batch Normalization vs. Instance Normalization



$$BN_f(x) = \gamma \frac{x - \mu_f(x)}{\sigma_f(x)} + \beta$$

channel-wise

$$IN_{n,f}(x) = \gamma \frac{x - \mu_{n,f}(x)}{\sigma_{n,f}(x)} + \beta$$

(sample,channel)-wise

