



K J Somaiya Institute of Technology

(Formerly known as K J Somaiya Institute of Engineering and Information Technology) An Autonomous Institute Permanently Affiliated to University of Mumbai.

Department of COMP Engineering

AI IA 2: Neural Style Transfer

Peeth Chowdhary - 16010122034 Rohit Deshpande - 16010122041 Eeshanya Joshi - 16010122074 Shubh Jalui - 16010122072 Ninad Marathe - 16010122106



Abstract

Neural Style Transfer (NST) merges the content of one image with the artistic style of another.

This project implements NST using two approaches: PyTorch-based optimization and TensorFlow Hub-based model inference.

Stylized outputs are generated using user-supplied images.



Problem Statement

Implement Neural Style Transfer for combining a content image with a style image

Achieve aesthetically pleasing style transfers using accessible tools and frameworks.

Objectives

- Implement NST using both optimization-based and pre-trained model-based techniques.
- Provide a user-friendly interface for uploading and processing images.
- Compare effects of different loss functions and hyperparameters on stylization quality.



Literature Survey

PUBLICATION YEAR	TITLE OF THE PAPER	OBJECTIVE
		To provide an overview
		of NSTs and the research
2017	Neural Style Transfer: A Review	conducted around them
		To test a different
		methodology of using
		VGG19 in ways that
		allow for flexible
	Dynamic Neural Style Transfer for Artistic Image	adjustments to style
	Generation using VGG19	weight ratios and reduce
2025		processing times



Implementation plan

PyTorch Version:

- Use a pre-trained VGG19 to extract content and style features.
- Compute style using Gram matrices; optimize a generated image to minimize style/content loss.

TensorFlow Version:

- Use TensorFlow Hub's pre-trained NST model for fast, high-quality transfer.
- Deploy with a simple Streamlit web interface for user interaction.

Tech Stack:

Languages: Python

Frameworks: PyTorch, TensorFlow, Streamlit

Libraries: OpenCV, PIL, NumPy, torchvision, TensorFlow Hub



```
import streamlit as st
import tensorflow as tf
import tensorflow hub as hub
import numpy as np
from PIL import Image
from io import BytesIO
# Set page config first
st.set_page_config(page_title="Neural Style Transfer", layout="wide")
# Cache the model
@st.cache resource
def load model():
    return hub.load('https://tfhub.dev/google/magenta/arbitrary-image-stylization-v1-256/2')
stylize model = load_model()
def load_image(image_file, image_size=(512, 512)):
    """Load and preprocess an image."""
    try:
       img = Image.open(image_file).convert('RGB')
       img = np.array(img, dtype=np.float32) / 255.0
        img = tf.image.resize(img, image_size, preserve_aspect_ratio=True)
       return img[tf.newaxis, :]
    except Exception as e:
        st.error(f"Error loading image: {e}")
       return None
def export_image(tf_img):
    """Convert tensor to downloadable PNG."""
    img = np.squeeze(tf img, axis=0) * 255
```

Deployment using
Streamlit for
real-time stylization
using
arbitrary-image-styli
zation-v1



```
def export_image(tf_img):
    """Convert tensor to downloadable PNG."""
   img = np.squeeze(tf img, axis=0) * 255
   pil_img = Image.fromarray(img.astype(np.uint8))
   buffer = BytesIO()
   pil img.save(buffer, format="PNG")
   return buffer.getvalue()
def main():
    """Streamlit UI for style transfer."""
    st.sidebar.title("Style Transfer")
    st.sidebar.write("Blend content with style!")
   # File uploaders
   content_file = st.sidebar.file_uploader("Content Image", ["jpg", "jpeg", "png"], help="Image to stylize")
    style_file = st.sidebar.file_uploader("Style Image", ["jpg", "jpeg", "png"], help="Style source")
    # Layout columns
    col1, col2, col3 = st.columns(3)
   # Process and display images
    content img = load image(content file) if content file else None
   if content file:
        col1.header("Content")
        col1.image(content file, use container width=True)
    style_img = load_image(style_file, (256, 256)) if style_file else None
    if style_file:
       col2.header("Style")
        col2.image(style file, use container width=True)
```



```
style_img = load_image(style_file, (256, 256)) if style_file else None
    if style file:
       col2.header("Style")
       col2.image(style_file, use_container_width=True)
   # Styling button
   if st.sidebar.button("Stylize"):
       if content_img is not None and style_img is not None:
           with st.spinner("Generating..."):
                stylized_img = stylize_model(content_img, style_img)[0].numpy()
                col3.header("Result")
                col3.image(stylized_img, use_container_width=True)
                col3.download button(
                    "Download",
                    export_image(stylized_img),
                   "stylized_image.png",
                    "image/png"
       else:
            st.sidebar.error("Upload both images.")
if name == " main ":
   main()
```



```
import os
import time
import torch
import torch.nn as nn
import torch.optim as optim
from PIL import Image
import torchvision.transforms as transforms
import torchvision.models as models
from torchvision.utils import save image
class VGG(nn.Module):
   def init (self):
       super(VGG, self). init ()
       self.chosen features = ['0', '5', '10', '19', '28']
       self.model = models.vgg19(pretrained=True).features[:29]
   def forward(self, x):
        features = []
        for layer num, layer in enumerate(self.model):
           x = layer(x)
           if str(layer num) in self.chosen features:
                features.append(x)
        return features
def load image(image name):
   image = Image.open(image name)
   image = loader(image).unsqueeze(0)
   return image.to(device)
def data collection(lr, m, fname):
   model = VGG().to(device).eval() # .eval() freezes the weights
   original image = load image('times square.jpg')
   style image = load image('style2.jpg')
```

Loss Function = Content loss (MSE or MAE) + Style loss using Gram matrices

Training Parameter = Varying learning rates (0.1, 1, 10)



```
def data_collection(lr, m, fname):
    model = VGG().to(device).eval() # .eval() freezes the weights
    original_image = load_image('times_square.jpg')
    style_image = load_image('style2.jpg')
    # generated_image = torch.randn(original_image.shape, device=device, requires_grad=True) # --> This is just noise
    generated_image = original_image.clone().requires_grad_(True) # Seems to work better than starting off as noise

# hyper parameters
    total_steps = 2000 # try 3000 and lower if taking too long
    learning_rate = lr # [0.001]
    alpha = 1 # content multiplier
    beta = 200 # style multiplier
    optimizer = optim.Adam([generated_image], lr=learning_rate)

file = open(fname+".csv", "w")
    line = 'step_number,total_loss,content_loss,style_loss,time'
    file.writelines(str(line) + '\n')
```



```
for step in range(total steps + 1):
   start = time.time()
   generated image features = model(generated image)
   original image features = model(original image)
   style image features = model(style image)
   style loss = content loss = 0
    for gen feature, orig feature, style feature in zip(
            generated image features, original image features, style image features):
       batch size, channel, height, width = gen feature.shape
       # Compute Gram Matrix for generated image
       gen matrix = gen feature.view(channel, height * width)
       gen gram matrix = gen matrix.mm(gen matrix.t())
       # Compute Gram Matrix for style image
       style matrix = style feature.view(channel, height * width)
       style gram matrix = style matrix.mm(style matrix.t())
       if m == 'rms':
            content loss = content loss + torch.mean((gen feature - orig feature) ** 2)
           style loss = style loss + torch.mean((gen gram matrix - style gram matrix) ** 2)
        else:
           mae loss = torch.nn.L1Loss()
           content loss = content loss + mae loss(gen feature, orig feature)
            style loss = style loss + mae loss(gen gram matrix, style gram matrix)
   total loss = alpha * content loss + beta * style loss
   optimizer.zero grad()
   total loss.backward()
   optimizer.step()
   t = f'Time: {time.time() - start}'
    if step % 50 == 0:
       print(total loss)
       save image(generated image, fname + '/' + fname + ',' + str(step) + '.png')
   line = str(step) + ',' + str(total loss.item()) + ',' + str(content loss.item()) + ',' + str(style loss.item()) + ',' + str(t)
   file.writelines(str(line) + '\n')
file.close()
```



```
name == ' main ':
device = torch.device('cpu' if not (torch.cuda.is available()) else 'cuda')
image size = 356 # use 178 = 356/2 or 89 if computation taking too long
loader = transforms.Compose(
       transforms.Resize((image size, image_size)),
       transforms.ToTensor(),
learning rates = [0.001, 0.01, 0.1]
metrics = ['mae', 'rms']
for i in range(1, 4):
    for lr in learning rates:
        for metric in metrics:
            file name = '{} {} t{}'.format(lr, metric, i)
           os.mkdir(file name)
           data collection(lr, metric, file name)
```



```
name == ' main ':
device = torch.device('cpu' if not (torch.cuda.is available()) else 'cuda')
image size = 356 # use 178 = 356/2 or 89 if computation taking too long
loader = transforms.Compose(
        transforms.Resize((image size, image size)),
        transforms.ToTensor(),
content multipliers = [0.001, 0.01, 0.1, 1, 10]
style multipliers = [x \text{ for } x \text{ in range}(200, 601, 100)]
for content multiplier in content multipliers:
    for style multiplier in style multipliers:
        file name = '{},{}'.format(content multiplier,style multiplier)
        os.mkdir(file name)
        data collection(content multiplier, style multiplier, file name)
```

Slightly altered code that changes the content and style multipliers (alpha + beta) instead of learning rates and ways to calculate loss.

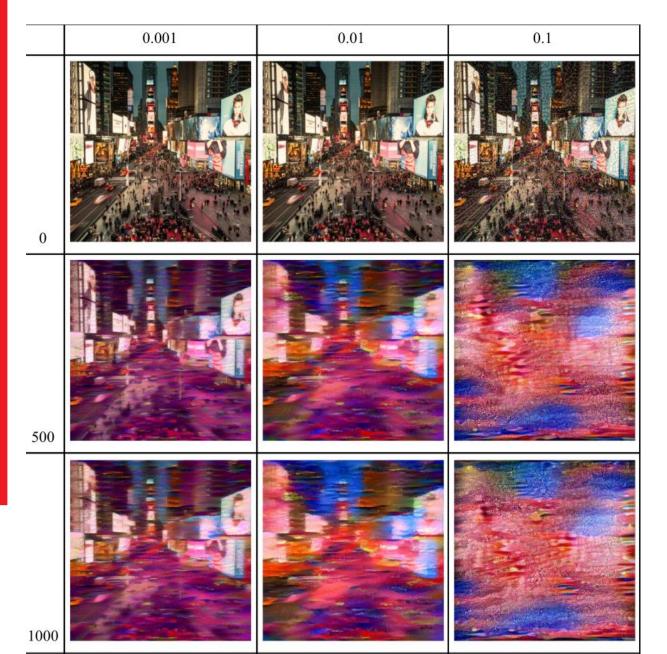


```
def data collection(a, b, fname):
    model = VGG().to(device).eval() # .eval() freezes the weights
    original image = load image('dog.jpeg')
    style image = load image('style.jpeg')
    # generated image = torch.randn(original image.shape, device=de
    generated image = original image.clone().requires grad (True)
    # hyper narameters
    total steps = 6000 # try 3000 a d lower if taking too long
    learning rate = 0.01 \# [0.001]
    alpha = a # content multiplier
    beta = b # style multiplier
    optimizer = optim.Adam([generated image], lr=learning rate)
    file = open(fname+".csv", "w")
    line = 'step number,total loss,content loss,style loss'
    file.writelines(str(line) + '\n')
```

Associated changes made to the beginning of the data_collection method



Implementation 2 Results



Changes in learning rate and how that affected the generated image over 100s of iterations



Conclusion

- Both optimization-based and model-inference approaches to NST are effective.
- The project showcases the flexibility and creativity enabled by neural networks in digital art generation.



References

• Studies mentioned in literature survey on slide 5