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Problem 1

1. Describe your implementation details and the difficulties you encountered.

按照 Pseudo code 實現 training 的過程:

```
def forward(self, x, c):
    """
    this method is used in training, so samples t and noise randomly
    """

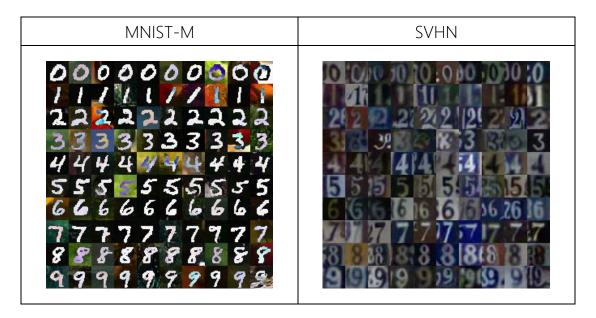
# t ~ Uniform(0, n_T)
    _ts = torch.randint(1, self.n_T+1, (x.shape[0],)).to(self.device)
# eps ~ N(0, 1)
noise = torch.randn_like(x)
# This is the x_t, which is sqrt(alphabar) x_0 + sqrt(1-alphabar) * eps
# We should predict the "error term" from this x_t. Loss is what we return.
x_t = (
        self.sqrtab[_ts, None, None, None] * x
        + self.sqrtmab[_ts, None, None, None] * noise
)
# dropout context with some probability
context_mask = torch.bernoulli(torch.zeros_like(c)+self.drop_prob).to(self.device)
# return MSE between added noise, and our predicted noise
return self.loss_mse(noise, self.nn_model(x_t, c, _ts / self.n_T, context_mask))
```

按照 Pseudo code 實現 Sampling 的過程:

```
# The reverse generation (denoising) process
for i in range(self.n_T, 0, -1):
   print(f'sampling timestep {i}',end='\r')
   t_is = torch.tensor([i / self.n_T]).to(device)
   t_is = t_is.repeat(n_sample,1,1,1)
   x_i = x_i.repeat(2,1,1,1)
   t_is = t_is.repeat(2,1,1,1)
   z = torch.randn(n_sample, *size).to(device) if i > 1 else 0
   # split predictions and compute weighting
   eps = self.nn_model(x_i, c_i, t_is, context_mask)
   eps1 = eps[:n_sample] # with condition
   eps2 = eps[n_sample:] # without condition
   eps = (1+guide_w)*eps1 - guide_w*eps2
   x_i = x_i[:n_sample]
   x_i = (
        self.oneover_sqrta[i] * (x_i - eps * self.mab_over_sqrtmab[i])
       + self.sqrt_beta_t[i] * z
   if i%80==0 or i==1:
       x_i_store.append(x_i.detach().cpu().numpy())
x_i_store = np.array(x_i_store)
return x_i, x_i_store
```

因為有找到可以參考的程式碼·所以在 implementation 上沒有遇到什麼大困難。 反而是因為運算資源不太夠所以 train 很久(最後還是屈服了·買了 Colab Pro)。

2. Please show 10 generated images for each digit (0-9) from both MNIST-M & SVHN dataset in your report. You can put all 100 outputs in one image with columns indicating different noise inputs and rows indicating different digits. [Visualize BOTH MNIST-M & SVHN]

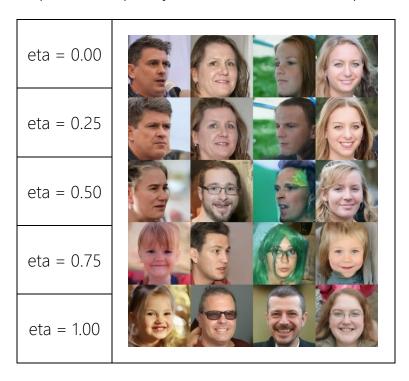


3. Visualize a total of six images from both MNIST-M & SVHN datasets in the reverse process of the first "0" in your outputs in (2) and with different time steps. [Visualize BOTH MNIST-M & SVHN]

MNIST-M					
t = 0	t = 80	t = 160	t = 240	t = 320	t = 400
				0	0
SVHN					
t = 300	t = 320	t = 340	t = 360	t = 380	t = 400
				10	10

Problem 2

1. Please generate face images of noise 00.pt ~ 03.pt with different eta in one grid. Report and explain your observation in this experiment.

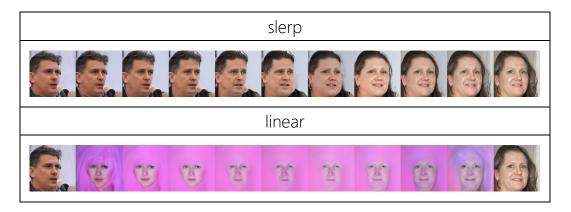


當 eta 越大的時候,加入的隨機 noise 就會對原本的 noise 影響越大。造成結果與 eta=0 的時候相差越遠,而每次 eta=0 所產生的的照片,因為沒有隨機性的關係,結果都會一樣。

2. Please generate the face images of the interpolation of noise 00.pt ~ 01.pt.

The interpolation formula is spherical linear interpolation, which is also known as slerp. What will happen if we simply use linear interpolation?

Explain and report your observation.



比起使用 slerp·linear 因為插值的過程中並沒有經過平滑·導致結果看起來會產生許多不自然的混和。而 slerp 則很明顯的正常許多。

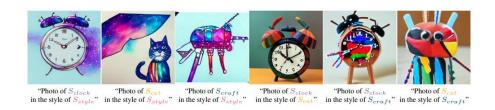
Problem 3

1. Conduct the CLIP-based zero shot classification on the hw2_data/clip_zeroshot/val, explain how CLIP do this, report the accuracy and 5 successful/failed cases.

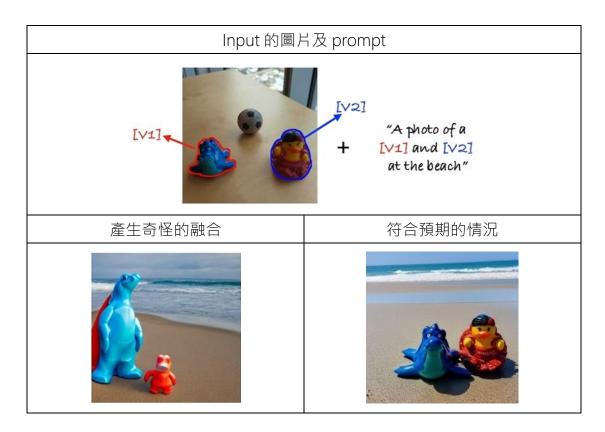
CLIP 如何做到 CLIP-based zero shot classification:

要對新的圖像進行分類時,CLIP 不需要針對目標類別重新訓練。而是可以直接為每個目標類別撰寫簡單的描述,像是「a photo of a cat」、「a photo of a dog」等等。然後,CLIP 將這些描述和待分類的圖像一起輸入模型,將它們都轉換成向量,並計算圖像向量和每個描述向量的相似度。模型最終會選擇相似度最高的描述,作為圖像的預測類別。

2. What will happen if you simply generate an image containing multiple concepts (e.g., a <new1> next to a <new2>)? You can use your own objects or the provided cat images in the dataset. Share your findings and survey a related paper that works on multiple concepts personalization, and share their method.

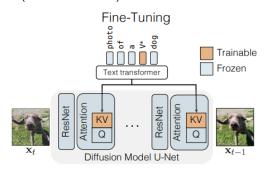


原 paper 有提到多個 concepts 分開放置的話是 OK 的,但如果 side-by-side 的放置的則不行。我有找到其他 paper 有呈現可能會產生一些奇怪的融合的情形,像下圖所展示的這樣。



Related paper that works on multiple concepts personalization:

我找到的是 Multi-Concept Customization of Text-to-Image Diffusion 這篇,這篇只對目標文本以及 Diffusion Model U-Net 的 cross-attention 層的 key 和 value 進行權重更新(如下圖所示)。



作法是將每一個 concept 分開訓練後再將各個概念的 cross-attention 投影矩陣 (key 和 value) 進行合併,合併的方式是使用約束最小平方問題(如下列式子·N是 concept 的數量·c 是攤平後的文本特徵)來讓不同的 concept 可以在互相影響最小的情況下合併,確保合併後仍能保有特徵。

$$\begin{split} \hat{W} &= \underset{W}{\operatorname{arg\,min}} ||WC_{\operatorname{reg}}^{\top} - W_0C_{\operatorname{reg}}^{\top}||_F \\ \text{s.t.} & WC^{\top} = V, \text{ where } C = [\mathbf{c}_1 \cdots \mathbf{c}_N]^{\top} \\ & \text{and } V = [W_1\mathbf{c}_1^{\top} \cdots W_N\mathbf{c}_N^{\top}]^{\top}. \end{split}$$

最後呈現的結果如下:



Reference

Problem 1: https://github.com/TeaPearce/Conditional_Diffusion_MNIST

Problem 2:

https://github.com/xiaohu2015/nngen/blob/main/models/diffusion_models/ddim_mnist.ipynb