Stable Diffusion

CLIP is used for text encoding.

CLIP algorithm explanation

```
# image_encoder - ResNet or Vision Transformer
# text_encoder - CBOW or Text Transformer
# I[n, h, w, c] - minibatch of aligned images
# T[n, 1] - minibatch of aligned texts
# W_i[d_i, d_e] - learned proj of image to embed
# W_t[d_t, d_e] - learned proj of text to embed
               - learned temperature parameter
# t
# extract feature representations of each modality
I_f = image_encoder(I) #[n, d_i]
T_f = text_encoder(T) #[n, d_t]
# joint multimodal embedding [n, d_e]
I_e = 12_normalize(np.dot(I_f, W_i), axis=1)
T_e = 12_normalize(np.dot(T_f, W_t), axis=1)
# scaled pairwise cosine similarities [n, n]
logits = np.dot(I_e, T_e.T) * np.exp(t)
# symmetric loss function
labels = np.arange(n)
loss_i = cross_entropy_loss(logits, labels, axis=0)
loss_t = cross_entropy_loss(logits, labels, axis=1)
loss = (loss_i + loss_t)/2
```

Figure 3. Numpy-like pseudocode for the core of an implementation of CLIP.

This pseudocode is a simplified core of **CLIP** (Contrastive Language–Image Pretraining), where the goal is to **align images and texts** in a shared embedding space. At a high level:

- Images and Texts get encoded separately into feature vectors.
- These feature vectors are projected into a **joint embedding space**.
- Then, they are **compared** using **cosine similarity**.
- A **contrastive loss** is applied, trying to **pull together** matching image-text pairs and **push apart** mismatched ones.

Now, let's go step-by-step:

1. Inputs

- I[n, h, w, c]: a minibatch of **n images** of shape height × width × channels.
- T[n, 1]: a minibatch of **n texts** (aligned, meaning *i-th image matches i-th text*).
- W_i[d_i, d_e]: a learned linear projection from image features to embedding space.
- W_t[d_t, d_e]: a learned linear projection from **text features** to **embedding space**.
- t : a learned temperature parameter (controls the sharpness of similarities).

2. Feature Extraction

```
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I_f = image_encoder(I) # [n, d_i]
T_f = text_encoder(T) # [n, d_t]
```

- The **image encoder** (like ResNet or ViT) outputs features of dimension d_i.
- The **text encoder** (like CBOW or Transformer) outputs features of dimension d_t.

3. Project into Joint Embedding Space

```
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I_e = 12_normalize(np.dot(I_f, W_i), axis=1)
T_e = 12_normalize(np.dot(T_f, W_t), axis=1)
```

- Multiply by learned projections W_i and W_t to map into a common space of dimension d_e.
- **L2 normalization** makes embeddings **unit vectors** (important because cosine similarity becomes simple dot product).

4. Compute Similarities

```
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logits = np.dot(I_e, T_e.T) * np.exp(t)
```

- Dot product between all **image embeddings** and all **text embeddings**.
- Result is a [n x n] similarity matrix: each entry [i, j] is the similarity between image i and text j.
- Multiply by exp(t), where t is learnable. Higher t → sharper softmax → harder contrastive learning.

5. Define the Labels

```
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labels = np.arange(n)
```

• This just says: for each sample *i*, the correct match is at position *i* (since images and texts are aligned).

6. Loss Computation

Here's the key part:

```
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loss_i = cross_entropy_loss(logits, labels, axis=0)
loss_t = cross_entropy_loss(logits, labels, axis=1)
loss = (loss_i + loss_t)/2
```

- loss_i:
 - For each **image**, treat its similarities to all texts as a softmax classification.
 - It should "predict" the correct text among all texts.
- loss_t:
 - o For each **text**, treat its similarities to all images as a softmax classification.
 - It should "predict" the correct image among all images.

▼ Take the average of the two: make sure both image-to-text and text-to-image alignments are good.

Why is this loss effective?

- It **pulls** corresponding image-text pairs closer together in the embedding space.
- It **pushes** non-matching pairs apart.
- It does **both directions**: *image* → *text* and *text* → *image*.
- Learning the **temperature** t allows it to dynamically adjust the difficulty.

Quick Intuition

If you have an image of a dog and a caption "a cute dog":

- The similarity between (dog image, "a cute dog") should be high.
- The similarity between (dog image, "a bowl of ramen") should be low.
- And vice-versa from text → image.

What is a Fixed Markov Chain

First, what is a Markov chain?

- A **Markov chain** is a sequence of random steps where:
 - The next state depends **only** on the **current** state (not the whole past).
- Mathematically:
 - P(xt+1|x0,x1,...,xt)=P(xt+1|xt)P(xt+1|x0,x1,...,xt)=P(xt+1|xt)
- This is called the **Markov property**: "memoryless" process.

You can imagine it as moving through a set of states step by step, according to some fixed transition probabilities.



Now, what is a fixed Markov chain with T steps?

Fixed:

• The transition rule (the probabilities of moving from one state to another) does not change over time.

T steps:

• You run the chain for **T time steps** (then stop).

In plain English:

A fixed Markov chain with T steps is a random process where you move between states according to the same fixed rules at every step, for exactly **T moves**.

At each time t=0,1,2,...,T-1t=0,1,2,...,T-1, you apply the same transition behavior, starting from some initial state.

After T steps, you stop.

Tiny Example

Imagine a 2-state Markov chain:

States:

- 0 = "Rainy"
- 1 = "Sunny"

Transition probabilities:

From / To Rainy (0) Sunny (1)

Rainy 0.7 0.3 Sunny 0.4 0.6

If it's rainy today, 70% chance rainy tomorrow.

• If it's sunny today, 60% chance sunny tomorrow.

Suppose T = 3.

You start Rainy (state 0):

- Step 1: pick next state based on Rainy row [0.7, 0.3]
- Step 2: pick next state based on new state's row
- Step 3: pick next state again

That's **T = 3 steps**, all using the same table.

The transition rule is fixed across steps.

Why does this matter?

- Many models (like diffusion models in machine learning) use fixed Markov chains to gradually corrupt data (forward process) or denoise it (reverse process).
- It's mathematically clean to reason about the evolution of distributions when transitions are fixed.
- If T is large, you can sometimes reach a **stationary distribution** (the system forgets its starting point).

Summary

Meaning Term

Markov chain Process where next state depends only on

current

Fixed Transition probabilities do not change over time

Prompt Integration with U-net

Let's go through **how Stable Diffusion uses cross-attention to integrate text into the U-Net** carefully.



In **Stable Diffusion**, the **U-Net** doesn't just denoise blindly.

Instead, at every step of denoising, it is guided by the text prompt.

This happens via **Cross-Attention** inside the U-Net.

Cross-Attention allows the U-Net to "look at" the text embeddings while deciding how to denoise.

@ Where exactly is Cross-Attention used?

- Inside the **middle and up-sampling blocks** of the U-Net.
- After certain convolution layers, there are self-attention blocks replaced or enhanced with cross-attentionblocks.

So the denoising is **conditioned** on the text **throughout the process**, not just at the start.



At a high level:

Componen t	What it is
Query (Q)	Comes from the current U-Net feature maps (the noisy image features)
Key (K)	Comes from the text embeddings (from the text encoder)
Value (V)	Comes from the text embeddings (same as keys)

The cross-attention computes:

Attention(Q,K,V) = softmax(QKTdk)VAttention(Q,K,V) = softmax(dkQKT)V

So, the image features "query" the text embeddings to decide what to pay attention to!

Why Cross-Attention Matters

Imagine you're trying to denoise an image towards "a cat wearing sunglasses":

- If the U-Net only looks at noise, it might not know which direction to denoise.
- But with cross-attention, it can **look at** the text "cat" and "sunglasses" words, **attend to them**, and guide the denoising **toward that meaning**.

Thus, cross-attention injects semantic meaning into the generation process at every level.



Suppose:

- Current U-Net feature map: $x \in \mathbb{R}^{n imes d_x}$ (n tokens, d_x dim)
- Text embeddings: $c \in \mathbb{R}^{m \times d_c}$ (m tokens, d_c dim)

Steps:

- 1. Project x into queries $Q = xW_Q$
- 2. Project c into keys and values: $K=cW_K, \quad V=cW_V$
- 3. Compute attention scores $A = \operatorname{softmax}(QK^T/\sqrt{d_k})$
- 4. Output: $A \times V$

This output is then added back into the U-Net feature maps, influencing the next steps of processing.

Result: Denoising is text-guided at every spatial location.

Tiny Flow Chart

Inside U-Net Block:

mathematica

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Noisy Feature Maps (Query) ---> Cross Attention ---> Updated Feature Maps
Text Embedding (Key & Value)

- U-Net feature maps act as queries.
- Text embeddings act as keys and values.
- Cross-attention lets the U-Net peek at the text and decide how to denoise better.
- This is **how text conditioning happens** in Stable Diffusion during image generation!

Stable Diffusion Use-cases and Architectural Adjustments

1. Mage Super-Resolution

(Make a blurry or low-res image sharper and higher-res)

Main idea:

Instead of starting from random noise, you condition the U-Net on a low-resolution image.

Architecture Changes:

- Input: Instead of pure noise xTxT, you start with a low-res noisy version of the image.
- Conditioning: Provide the low-resolution image as an extra input along with noise and timestep.
- Sometimes the low-res image is **concatenated** with the noisy latent as extra channels, or passed via cross-attention.
- The U-Net learns to "refine" the blurry image progressively.

In short:

Low-res image is treated like a guiding prompt, similar to how text was used.

2. 🎨 Style Transfer

(Apply the style of one image onto another)

Main idea:

You condition the diffusion model on two things:

- The **content image** (what you want to preserve)
- The **style image** (what artistic style you want to apply)

Architecture Changes:

- **Conditioning:** Both the content and style information are encoded.
 - Content image could be injected into U-Net feature maps.
 - Style information could influence cross-attention layers.
- Objective: Generate an image that matches content features of the content image and style features of the style image.

- Sometimes two encoders are used:
 - o One for content, one for style.

In short:

You combine **dual conditioning** inside the U-Net: content + style.

3. Inpainting

(Fill missing parts of an image, e.g., remove object and refill background)

Main idea:

You give the model:

- A partially masked image
- A mask that tells which regions are missing

Architecture Changes:

- Input: Instead of only random noise, the input includes:
 - Noisy image + mask
- Conditioning:
 - Masked regions are where model needs to "invent" data.
 - Unmasked regions (known parts) are copied from the input.
- **Training:** During training, you randomly mask patches in images so the model learns to inpaint missing regions.

Typical implementation:

- Concatenate the image + mask as input channels.
- Model predicts denoised latents conditioned on visible context.

In short:

Mask + visible parts tell the U-Net where to "hallucinate" new pixels.

4. **M** Outpainting

(Extend an image beyond its original boundaries)

Main idea:

Similar to inpainting, but **now you add new empty regions around the image**, and the model **hallucinates outward**.

Architecture Changes:

- **Input:** Image placed in the center of a larger canvas, with surrounding regions as masked (empty).
- Conditioning:
 - o Original image guides the newly generated parts.
 - o Empty surrounding regions are where the U-Net needs to generate new content.
- Sometimes the text prompt is also still used, to control what the extended scene should contain.

In short:

Exactly like inpainting, but the **mask** is applied around the original image instead of within.

Common Pattern Across All These Tasks

Aspect How it Changes

Input Not just random noise — it's noise + guiding context (text, images, masks)

Conditioning U-Net is modified to accept extra conditions, either via concatenation or

cross-attention

Training Needs slightly modified tasks (e.g., masked training for inpainting)

Loss Often still predict noise (like in basic diffusion), but may add extra perceptual

or style losses

▼ The U-Net remains the core engine,



Use Case How U-Net Differs

Super-Resolutio Condition on low-res image

Style Transfer Condition on content + style encodings

Inpainting Condition on masked image + mask

Outpainting Same as inpainting but mask surrounds the

image