

# DreamBooth Training Guide: Cinematic two Person Conversion

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- **Goal**

Train a Stable Diffusion DreamBooth model using the notebook to generate cinematic-style images of a specific person, suitable for character conversations or movie-like storytelling.

- **Prerequisites**

- Google Account
- Basic knowledge of Colab & Hugging Face
- A few high-quality images (3–10) of the target person
- Hugging Face token (for downloading base models)

- **Steps to Train:**

1. Open the Notebook:

Use the provided link to open in Google Colab.:

[https://colab.research.google.com/drive/1GW-j1\\_5nSM7YjByBgOdXtJMk9b\\_Yt2vv?authuser=1](https://colab.research.google.com/drive/1GW-j1_5nSM7YjByBgOdXtJMk9b_Yt2vv?authuser=1)

2. Setup Environment:

Run the initial cells to install dependencies (diffusers, accelerate, etc.).

3. Install Requirements:

```
!wget -q
https://gist.githubusercontent.com/FurkanGozukara/be7be5f9f7820d0bb85a3052
874f184e/raw/d8d179da6cab0735bd5832029c2dec5163db87b4/train_dreamboo
th.py
!wget -q
https://github.com/ShivamShrirao/diffusers/raw/main/scripts/convert_diffusers_
to_original_stable_diffusion.py
%pip uninstall torchtext --yes
%pip install -qq git+https://github.com/ShivamShrirao/diffusers
%pip install torch==2.2.0 torchvision torchaudio --index-url
https://download.pytorch.org/whl/cu121 --upgrade
%pip install -q -U --pre triton --upgrade
%pip install -q accelerate transformers ftfy gradio natsort safetensors
%pip install bitsandbytes==0.41.3 --upgrade
%pip install xformers==0.0.24 --upgrade
```

```
%pip install triton==2.2.0 --upgrade
%pip install diffusers==0.27.0 --upgrade
%pip install huggingface_hub==0.25.2
%pip install numpy==1.26.4
%pip install transformers==4.43
```

4. Connect drive, download stable diffusion model
5. Login to Hugging Face:  
Enter your Hugging Face token when prompted (sign up at <https://huggingface.co>).
6. Upload Your Images:  
Prepare 3–10 photos of the person in cinematic lighting (portrait, different angles).  
Upload to a folder named something like 'person\_cinematic/'.
7. Define Your Prompt Template:  
Use a custom token like 'cinematic [name]' or 'photo of cinematic man' in your training prompt. Example: 'a cinematic portrait of sks person'.
8. Set Training Parameters:
  - Pretrained model: runwayml/stable-diffusion-v1-5
  - Resolution: 512x512
  - Steps: ~800–2000 (depending on overfitting)
  - Learning Rate: ~5e-6
  - Save checkpoints every 500 steps
9. Run Training Cell:  
It will train the model using the uploaded images and prompt.
10. Run Gradio UI for generating images.
11. Test Generation:  
Use the inference cell to try prompts like: 'a cinematic close-up of sks person, dramatic lighting, 35mm film'
12. Optional) Delete diffuser and old weights and only keep the ckpt to free up drive space.

## • Tips

- Use cinematic tags: 'cinematic lighting', 'moody', 'film grain', 'anamorphic', 'bokeh'.
- Maintain 1:1 face-to-frame size for input images.
- Use prompt editing later for expressions or different camera angles.

## • Links

Hugging Face: <https://huggingface.co>

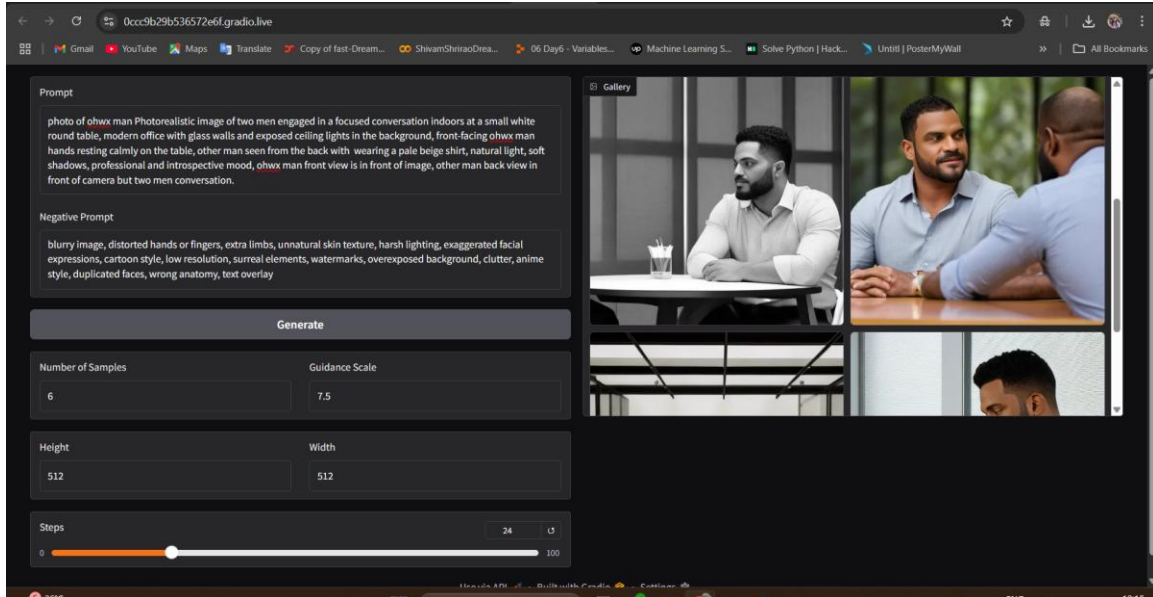
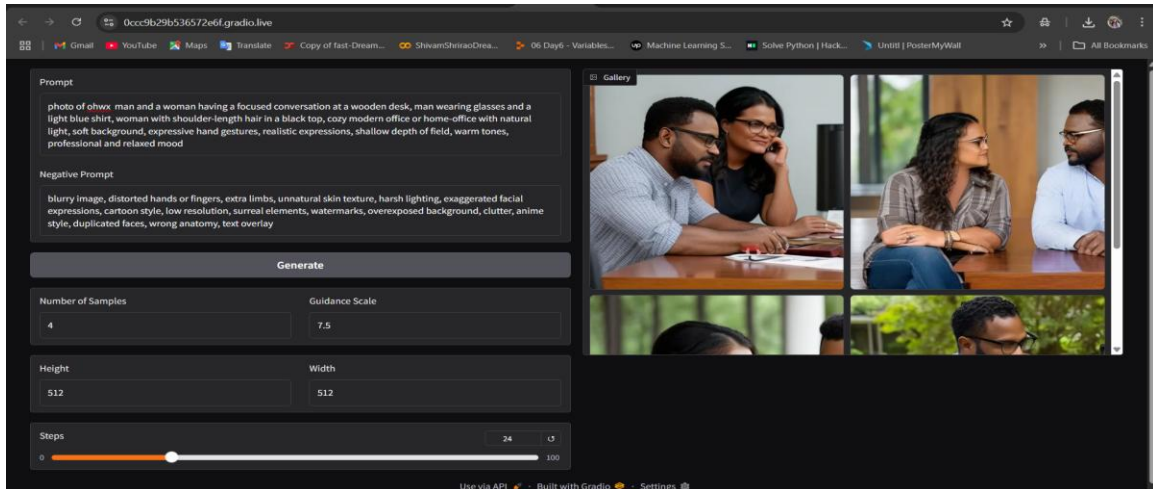
Stable Diffusion Model: stable-diffusion-v1-5/stable-diffusion-v1-5:  
<https://huggingface.co/stable-diffusion-v1-5/stable-diffusion-v1-5>

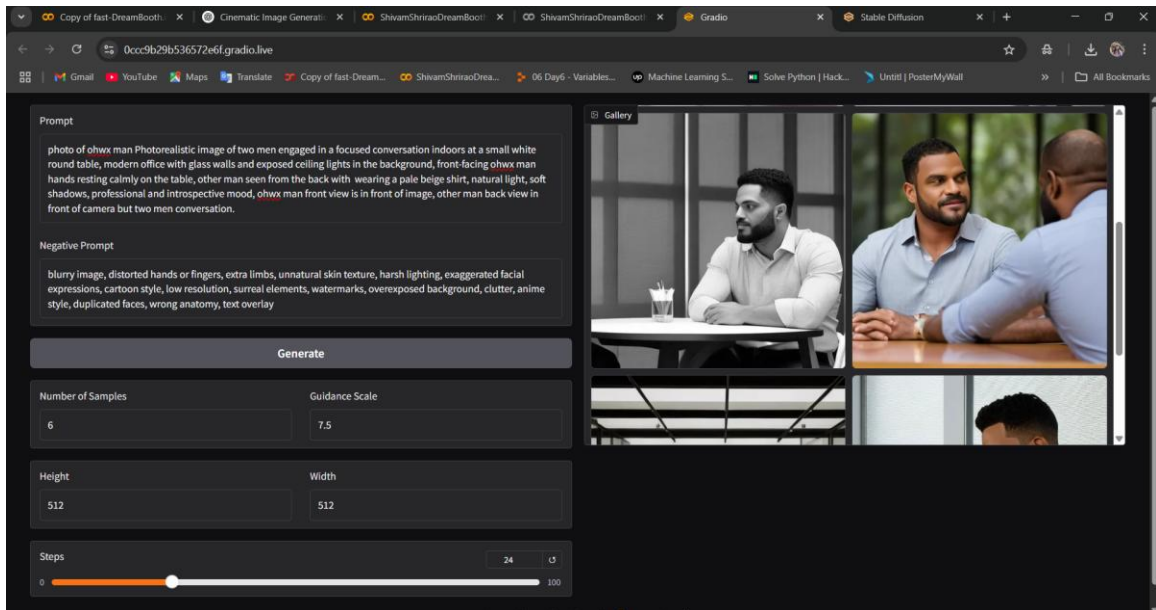
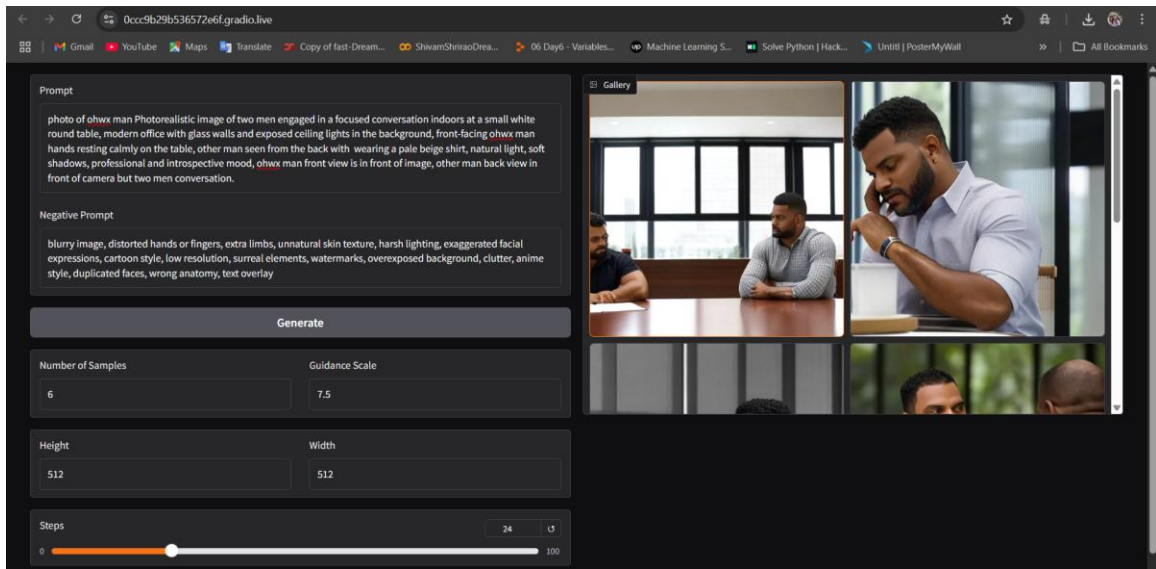
DreamBooth Colab Notebook:

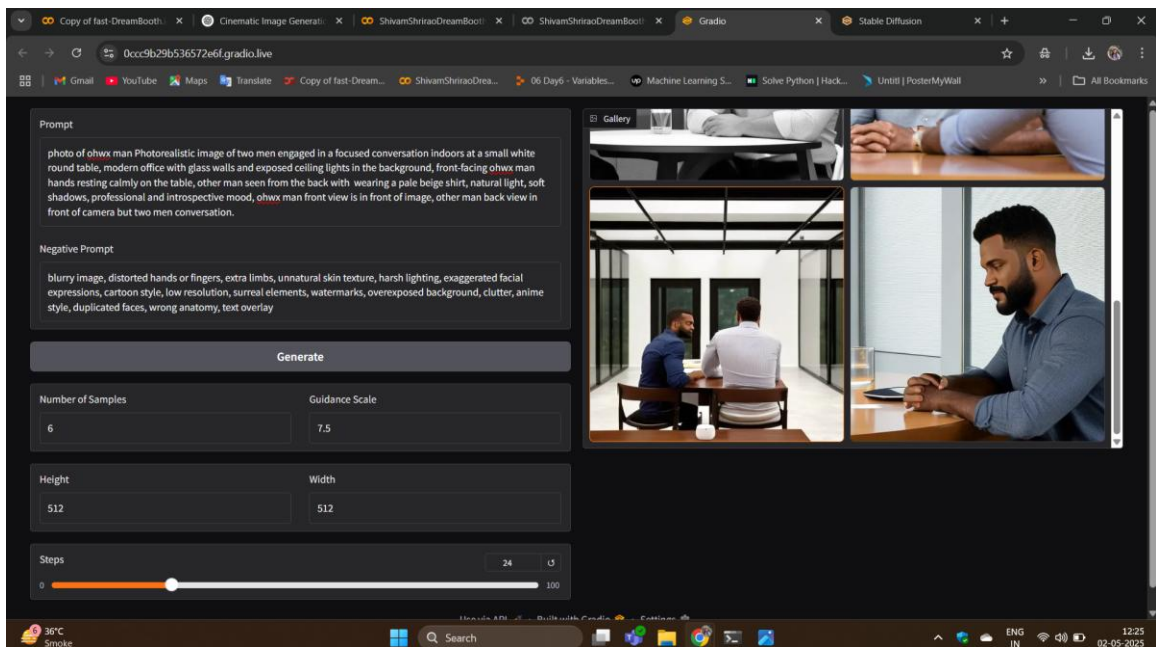
<https://colab.research.google.com/github/FurkanGozukara/Stable-Diffusion/blob/main/DreamBooth/ShivamShriraoDreamBooth.ipynb?authuser=1>

- **Outputs of model**

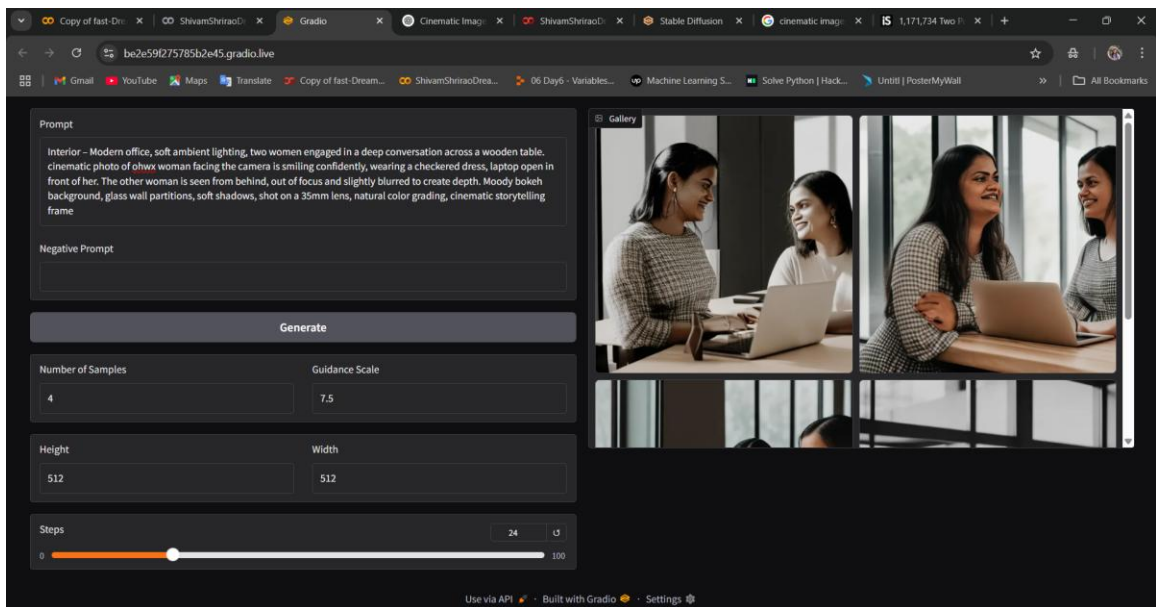
Ohwx man:







Ohwx woman:





**vrthk man:**



