

Stable Diffusion *Dream Studio* Guide



Introduction

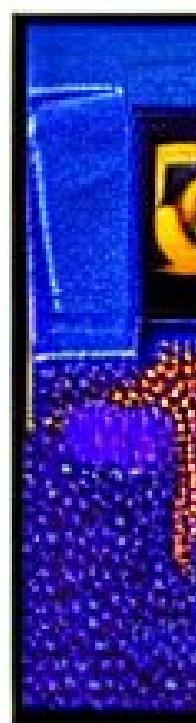
We will go over a variety of different subjects. Starting, of course, with the basics. This way you can understand the process a bit better if you're completely new to generating Stable Diffusion art through Dream Studio. Later on you will see some examples of different image sizes, CFG scales, prompts and more. To give you a little bit of an idea of what you can expect.

Some of the images in this guide have been created when Stable Diffusion still worked through the Discord in beta 1, and others have been made through the website Dream Studio. This has no difference, since both work with the same model. The Dream Studio is simply a new way of interacting with the service.

Please keep in mind that, even though we go over settings and prompts in this guide, there is no one answer that fits all. Many times you think you've found the perfect settings, e.g. CFG scale, steps, image size... But then you write a different prompt and all of a sudden those settings create very bad results. Even though this guide is meant to give you an idea of possibilities with Stable Diffusion, and to some extent some tips and tricks. Please know that we are all still learning, and it's very much a case of experimentation to see what works best for certain styles/prompts/outcomes.

(I am not affiliated with Stable Diffusion / Stability.ai, but simply a curious nerd who wanted to share his experience and possibly provide some useful info. I hope you enjoy it.

Happy creating! - You can reach out to me if you got any questions or suggestions, MrH3RB#9625 on discord, @MrH3RB on twitter, or @CaptainH3RB on reddit.)



Why Use Artificial Intelligence?

These tools are still quite new, and can even be a little scary at times. Will traditional artists lose their jobs? Is art even art when a computer makes it? What's the point even, type in some words and the machine spits out an image...

Well, I don't have all the answers. But yeh, it can be scary. No I don't believe artists will lose their jobs, if anything these new tools will bring the ability to adapt to new ways of creativity that haven't even been thought of before. And what is the point of all this? Creating with the machine? Well it's a lot more fun doing it, than it looks like when you just hear about these things. We mostly see text to image tools like Stable Diffusion, DALL-E 2 and MidJourney. But this is about so much more than just making images.

We are at a state in life where the ability to create is about to change massively. And that is just the beginning. These AI systems are capable of much more which we will all witness soon enough. But let's not get ahead of ourselves.

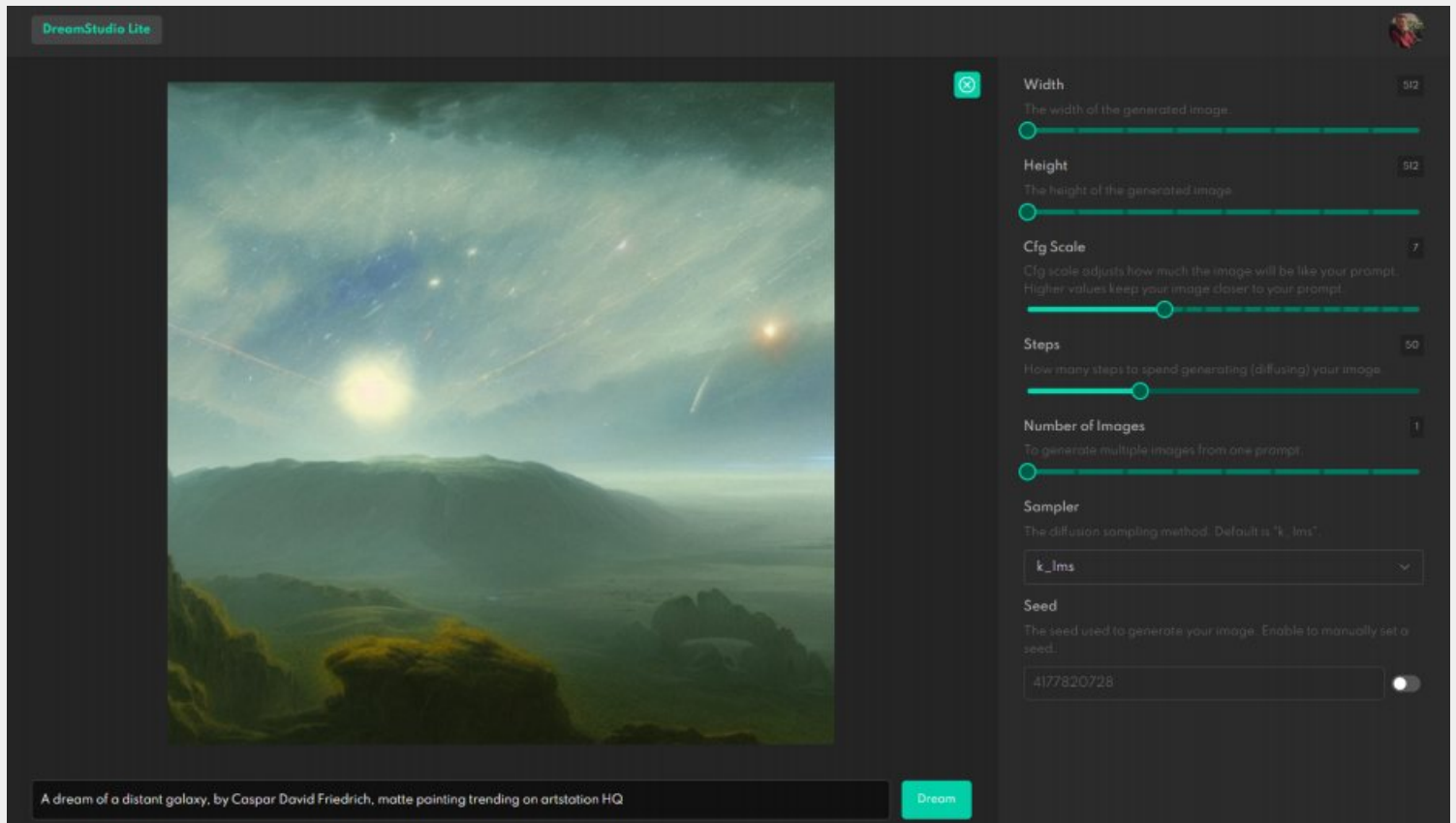
Not everybody is able to create in the traditional way. Whether that is drawing or painting, or making pottery or anything else. These new tools will make it so that millions and millions of people all of a sudden have the ability to create what they have in mind, without a high barrier of entry. Or maybe they just couldn't because of a physical condition, and now there is this new tool that brings the opportunity to still get their amazing thoughts into some form of creativity.

I'm not here to preach, but I did want to write this little part into this guide. Because I strongly feel that there is no need to compete. No need for either or, as in only traditional art or only ai art can prevail. Since I'm someone who tried to learn traditional art, but can't because of physical conditions I can somewhat see both sides of the arguments being made online. And I think we can all get along and create beautiful things together.

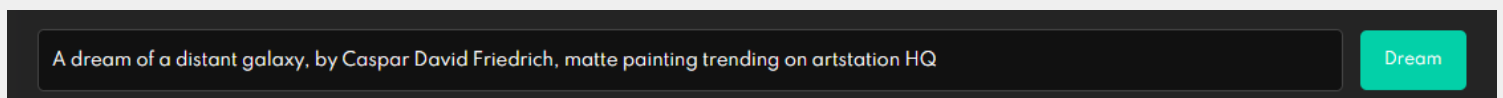
The Basics

Generating images in [Dream Studio](#) might seem complicated, with all the adjustable sliders. But it isn't actually that difficult.

Prompting your images



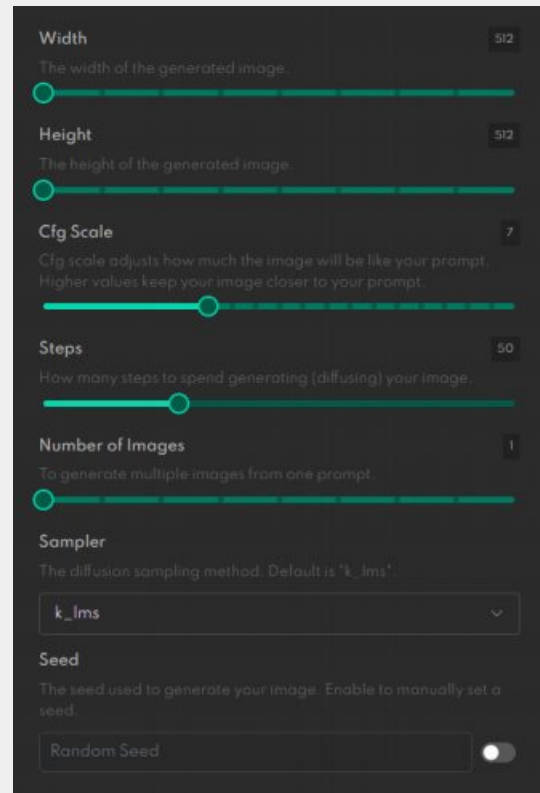
At the bottom of the screen you can fill in your prompt, after which you hit the “***Dream***” button, and Stable Diffusion will start generating your image.



Generating images through Dream Studio works with a credit system.

The default settings will cost you 1 credit per generated image.

Default settings:



The screenshot shows the Dream Studio settings panel with the following default values:

- Width:** 512 (The width of the generated image.)
- Height:** 512 (The height of the generated image.)
- CFG Scale:** 7 (CFG scale adjusts how much the image will be like your prompt. Higher values keep your image closer to your prompt.)
- Steps:** 50 (How many steps to spend generating (diffusing) your image.)
- Number of Images:** 1 (To generate multiple images from one prompt.)
- Sampler:** k_lms (The diffusion sampling method. Default is 'k_lms'.)
- Seed:** Random Seed (The seed used to generate your image. Enable to manually set a seed.)

Changing the amount of steps and / or the Width / Height of your image will cost you more credits. You can adjust the CFG, Sampler and seed without additional costs.

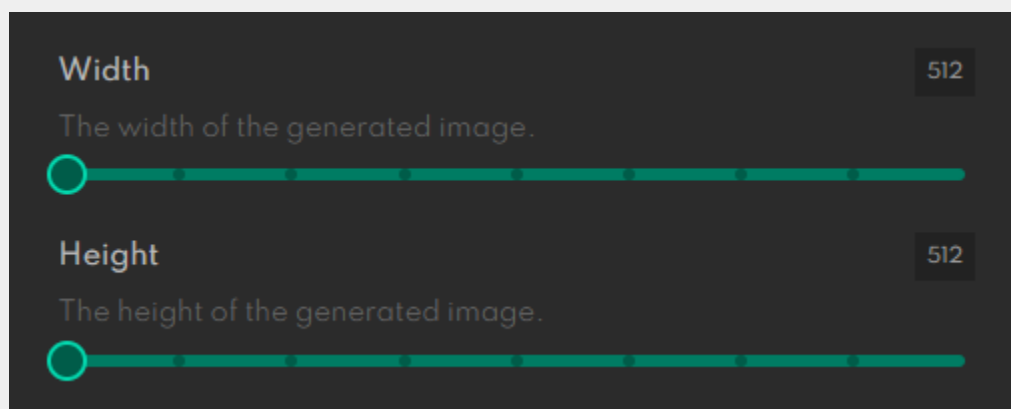
See the graph below for more information:

Compute Costs & Generation Counts Explained						
Steps	512x512	512x768	512x1024	768x768	768x1024	1024x1024
10	0.2	0.5	0.8	0.9	1.3	1.9
25	0.5	1.2	1.9	2.3	3.3	4.7
50	1.0	2.4	3.8	4.6	6.6	9.4
75	1.5	3.6	5.7	6.9	9.9	14.1
100	2.0	4.8	7.6	9.2	13.2	18.8
150	3.0	7.2	11.4	13.8	19.8	28.2

What are the different settings?

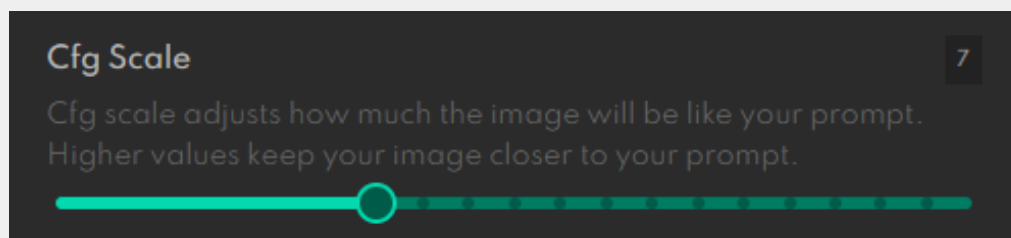
Let's go over the settings briefly. Later on we will look at some generated images with different settings to see how they affect the images in more detail.

Width & Height:



This one is pretty self explanatory, so we will move on to the next setting.

CFG Scale:



The CFG scale basically tells Stable Diffusion how closely it has to follow your prompt input. You can see some examples of this later on. As default this is set to seven (7), which is fine for most prompts. You can experiment with this, because with some prompts it might generate a much better result when you put the scale a bit higher. Do note however that increasing the CFG scale might cause artifacts in your images. This, again, does depend on your prompt. You can see examples of this later on.

Steps:

Steps 50

How many steps to spend generating (diffusing) your image.

Increasing the amount of steps tells Stable Diffusion that it should take more steps to generate your final result which can increase the amount of detail in your image. Higher steps does not always equal a better result. Especially when using a huge amount, like 100-150 or even higher. Just like the CFG scale, this does depend on your prompt, but going too high with the steps can cause artifacts in your final result.

Do note that increasing the amount of steps will cost you more credits per image. (see the graph on the bottom of page four (4).)

Number of Images:

Number of Images 1

To generate multiple images from one prompt.

Increasing the number of images will create more images with the same prompt you have typed.

Sampler:

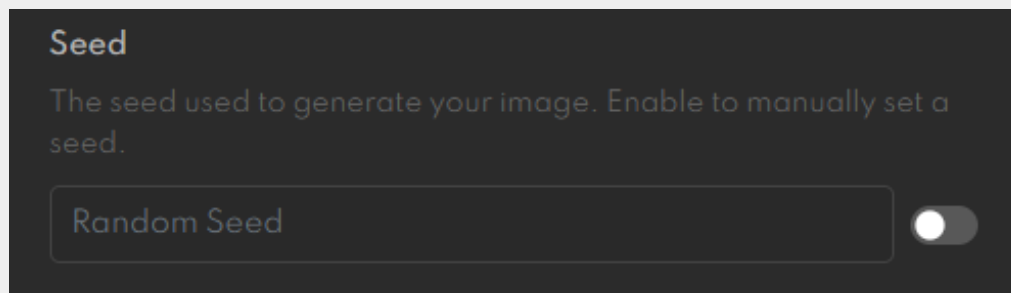
Sampler

The diffusion sampling method. Default is "k_lms".

k_lms

The sampler is what Stable Diffusion uses to decide how to generate your final result. You can change this, but the differences between samplers are often very small. But again, it depends a bit on your prompt.

Seed:



The seed determines the starting point of your image generation. Normally a seed is randomly chosen when you hit the “Dream” button. But if you find something you like, you can use the same seed to get the same image again. Do keep in mind that to get the exact same result as before you need to not only use the same seed, but also the same settings and prompt!

Something that is fun about this seed system, is that you can change some small things about your prompt while keeping the same seed, and thus generate very similar looking images. You’ll see some examples of this later on.

Width & Height Examples

You can experiment with the width/height as much as you want.

Just keep in mind that changing this will change your composition, and costs extra credits.



Another note on the image size:

The AI has been trained on countless 512x512 images. Causing it to sometimes create doubles when you try to generate an image that has a different size. Like a neck on top of a head, as the image here:



The explanation from the developers has been that, because the AI has been trained on images that were 512x512, it will try to generate your prompt in every 512x512 area. Meaning, if you ask for a portrait it will try to create a portrait in every square with the above mentioned dimensions. Whatever your prompt is, you might run into duplicates or odd images like the one above when you go above the default image size.

This does not always happen, but it is something to be aware of. Especially because bigger images cost more credits. *(see the graph on the bottom of page four (4).)*

Composition does change based on your image size, so if you have another composition in mind it could be good to experiment with other sizes than the default 512x512.

If you simply want to generate bigger images for quality purposes I'd advise against that given the extra cost. In that case it's better to just create with the 512x512 settings and upscale the image with a service like BigJPG, or a google colab notebook or whatever upscale service you prefer. Many of these upscale services use AI, and have barely any quality loss, if any.







































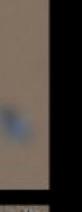














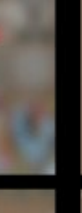
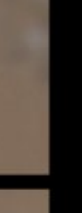






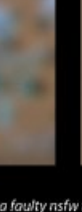
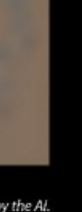




Stable Diffusion CFG /Sampler Comparison

!dream "Andromeda galaxy in a bottle" -H 640 -W 256 -C 1 -S 328034772

12

	k_lms	k_heun	k_euler	k_euler_ancestral	k_dpm_2	k_dpm_2_ancestral	ddim	plms
CFG 3								
CFG 5								
CFG 7								
CFG 10								
CFG 15								
CFG 20								
CFG 25								
CFG 30								

*Blurred images are because of a faulty nsfw flag by the AI.

Generation Steps Examples

It might seem logical to always run at the maximum amount of steps, but this isn't always a good idea. Often you don't see that much of a difference when running your steps higher than 70-100, depending on your prompts. And in most cases, for simple images, 50 is plenty. [Link](#) to image.

(CFG scale was set to 15)

Stable Diffusion Step Comparison

!dream "surrealism, beautiful face of a Russian princess. skull mask. dark eyes. intricate detail. god rays. octane render, trending on artstation. cinematic. hyper realism. high detail. 8K. iridescent details." -S 34449963 k_dpm_2_ancestral



Random vs Same Seed Examples

When you keep the seed random you will get random starting noise, resulting in different images. But if you want to adjust an image you like, you can keep the same seed as explained on page seven (7).

Here you see an example of the same prompt with four different outcomes. (Top left grid image).

When you keep changing your prompt you can slightly adjust images you like. (See the next images.)

Stable Diffusion Seed Comparison

!dream "glowing mushroom character with big oval shaped eyes. sharp focus. highly detailed. cute. vibrant. pixar." -C 35 -S 3892359291 k_dpm_2



Stable Diffusion Seed Comparison

!dream "glowing mushroom character with big oval shaped eyes.
sharp focus. highly detailed. cute. pixar." -C 25 -S 3892359291

k_dpm_2



I've only changed two things, yet the images are quite different. (Removed "vibrant" from the prompt, and lowered the CFG scale from 35 to 25.)

Experiment a bit with the seeds, and your prompts. You can get amazing results when you change things around a bit. A tip, don't change too much at once. The more variables you add, the harder it is to actually find out what created the image you so enjoy.

Change one thing, hit *Dream*, reflect on the changes, and repeat. That is how you learn!

I hope this guide helped you create better art, or at least understand the tool a bit more. Below are some other resources you can take a look at! If you have any questions, suggestions or want me to add something to this guide, let me know.

Go create something: beta.dreamstudio.ai

Other Resources

- Stable Diffusion [Akashic Records](#)
- Lighting prompts [comparison](#)
- Stable Diffusion [guide](#) for non-artists.
- Stable Diffusion ultimate [beginners guide](#).
- Dream Studio [tips](#): Maximize your credits.
- DALL-E 2 prompt [book](#)