

Exploiting the Signal-Leak Bias in Diffusion Models

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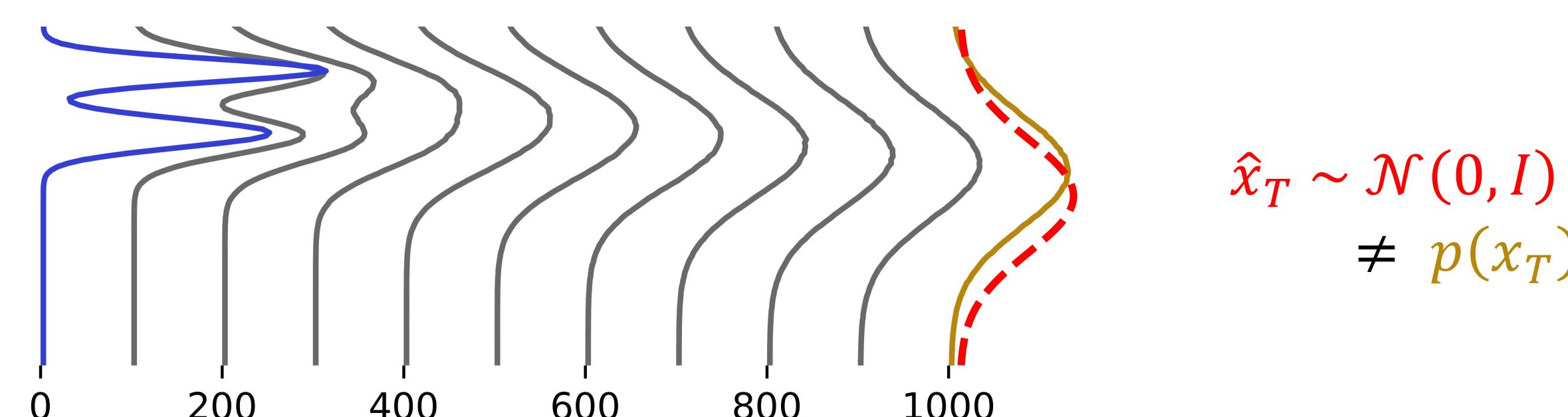
Signal-leak bias

We can generate images in a **desired style** or with a more natural color distribution **without retraining** the diffusion model, by exploiting a **signal-leak bias** present in the model.

Common **diffusion models** never fully corrupt **images** during training [1,2]:

$$x_T = \sqrt{\bar{\alpha}_T} x_0 + \sqrt{1 - \bar{\alpha}_T} \varepsilon \text{ with } x_0 \sim p(x_0) \text{ and } \varepsilon \sim \mathcal{N}(0, I)$$

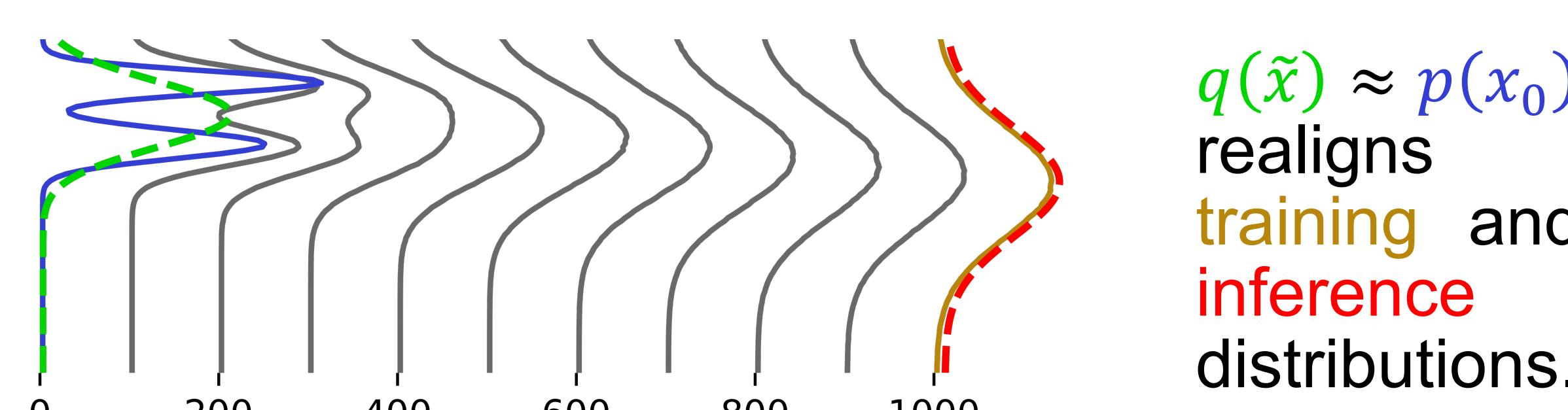
However, the process of **generating images** starts with **pure noise** $\hat{x}_T \sim \mathcal{N}(0, I)$, oblivious of the **signal leak** $\sqrt{\bar{\alpha}_T} x_0$ present in x_T during training, **creating a bias**.



Instead of retraining or finetuning [1,2,3] to remove this bias, we exploit it to our advantage, generating images in the style we want.

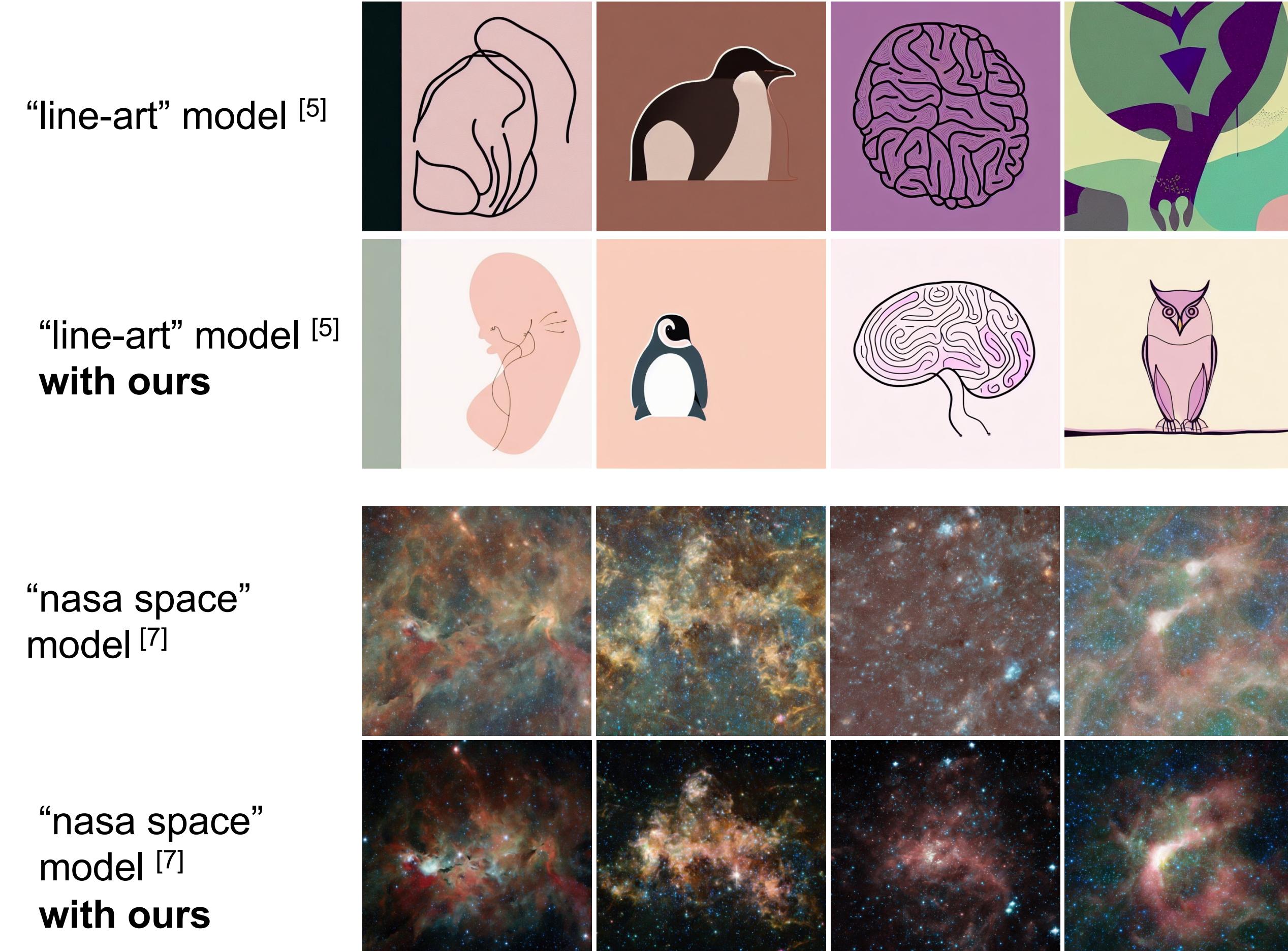
We include a **signal-leak** $\sqrt{\bar{\alpha}_T} \tilde{x}$ in \hat{x}_T at **inference time**, starting generating images from:

$$\hat{x}_T = \sqrt{\bar{\alpha}_T} \tilde{x} + \sqrt{1 - \bar{\alpha}_T} \varepsilon \text{ with } \tilde{x} \sim q(\tilde{x}) \text{ and } \varepsilon \sim \mathcal{N}(0, I)$$



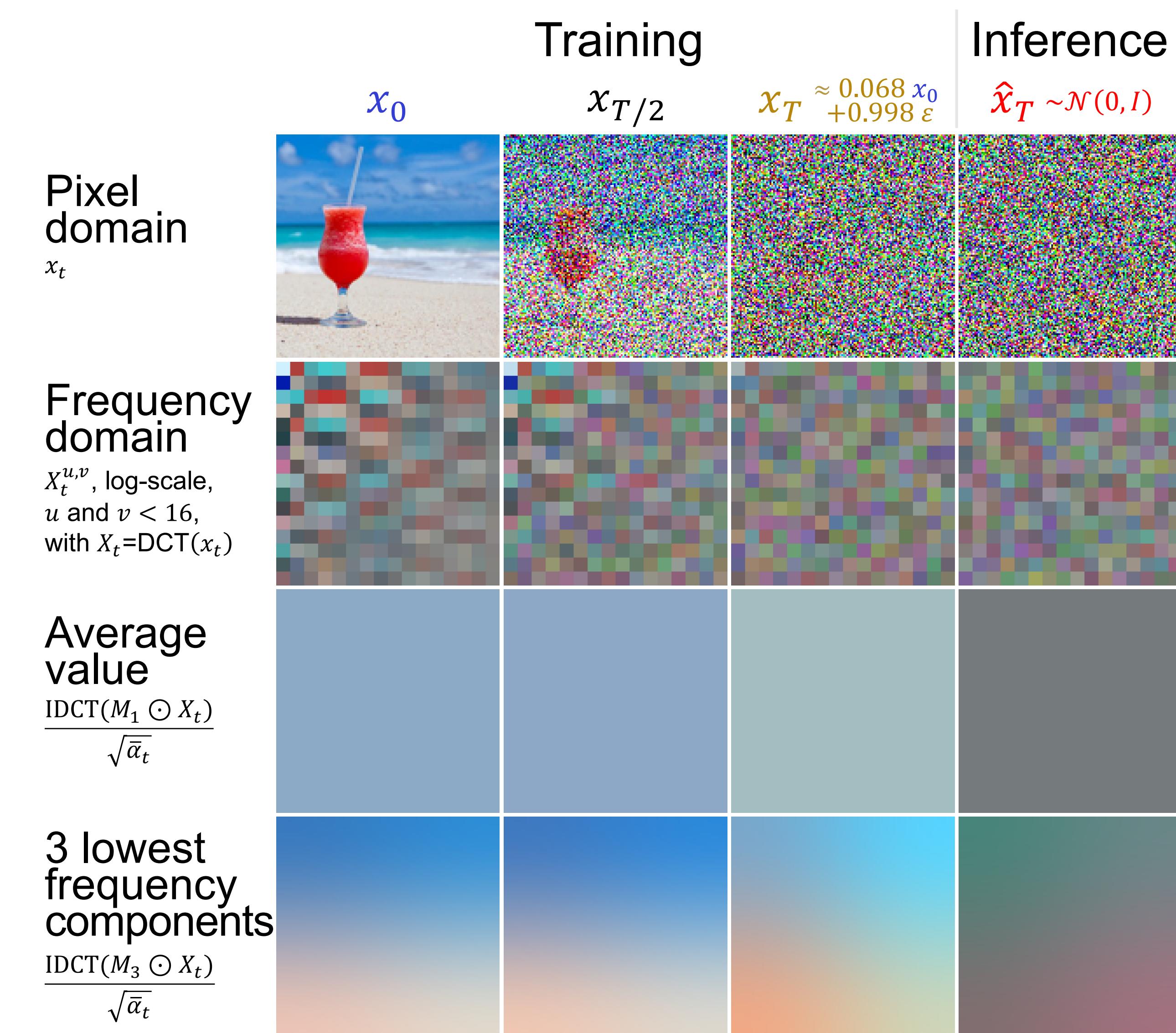
Fixing style-adapted models

We obtain a distribution $q(\tilde{x})$ in the **pixel domain**, by approximating the distribution $p(x_0)$ as independent Gaussian distributions for each pixel.



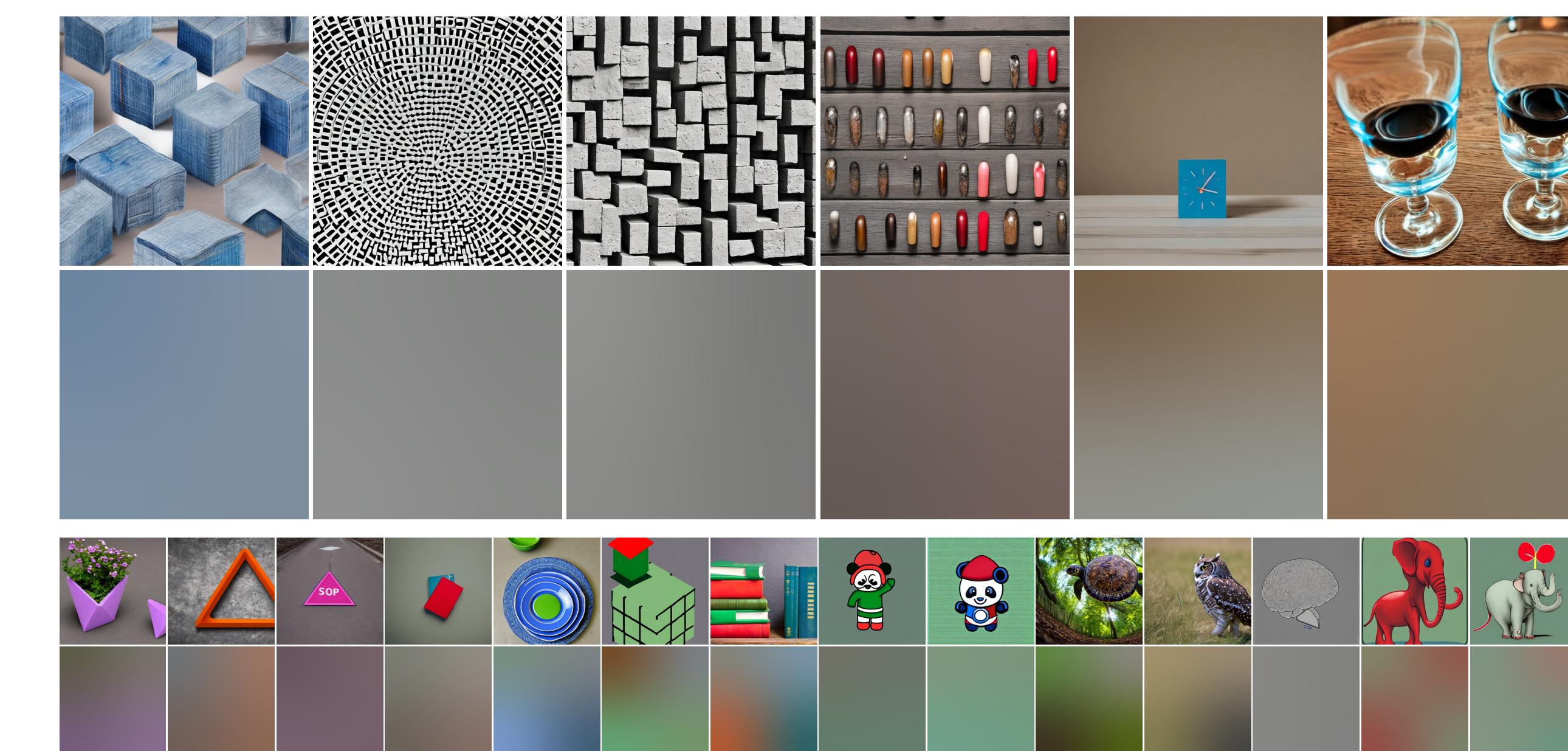
Better low-frequency components

The diffusion model uses the signal-leak $\sqrt{\bar{\alpha}_T} x_0$ to deduce the **low-frequency information** about x_0 from x_T . Using $\hat{x}_T \sim \mathcal{N}(0, I)$ biases the low-frequency components towards **medium values**.



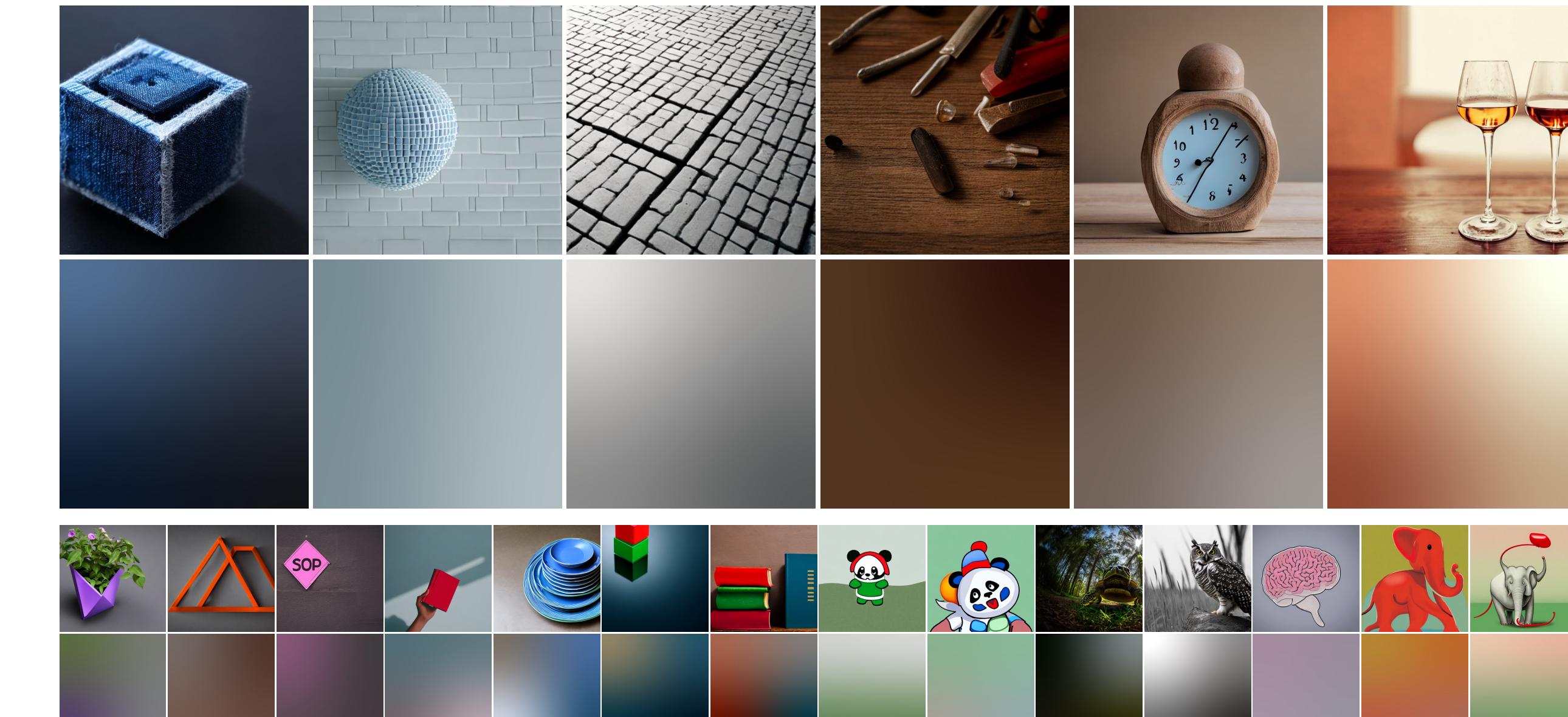
To avoid this, we additionally **model the low-frequency components**, estimating their mean and covariance, and obtain a distribution $q(\tilde{x}) \approx p(x_0)$.

Original results SD 2.1 [4] → greyish images with low contrast or variation of colors



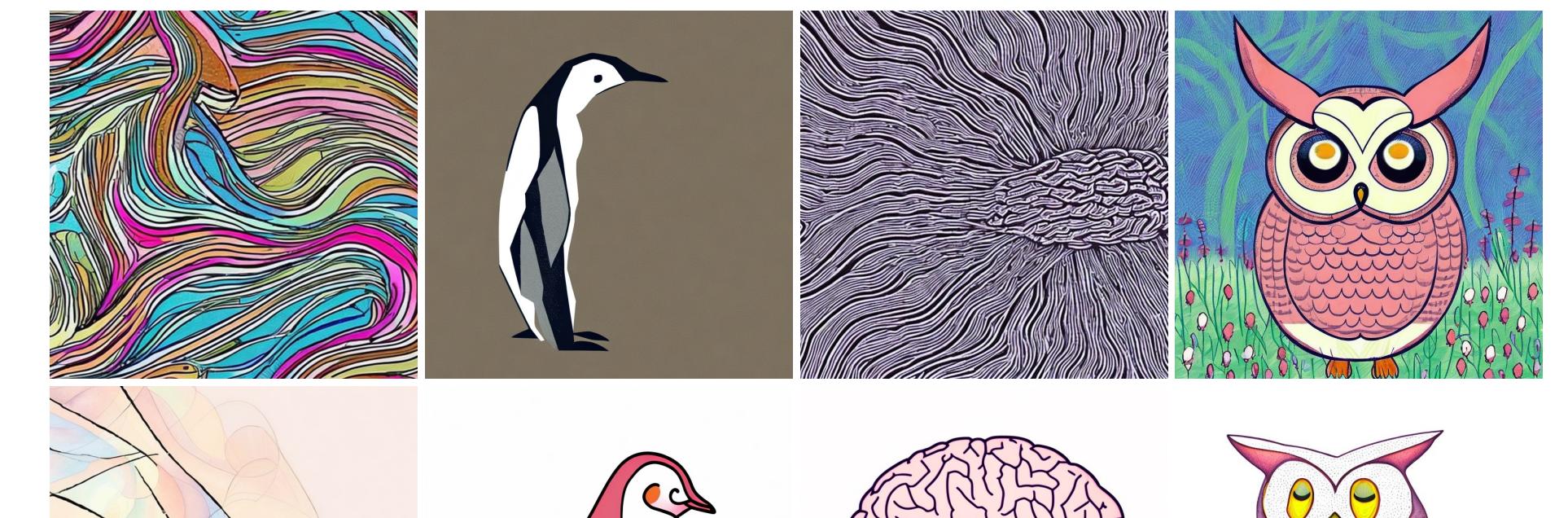
Our results

SD 2.1 [4] **with ours** → more varied and natural distribution of low-frequency components

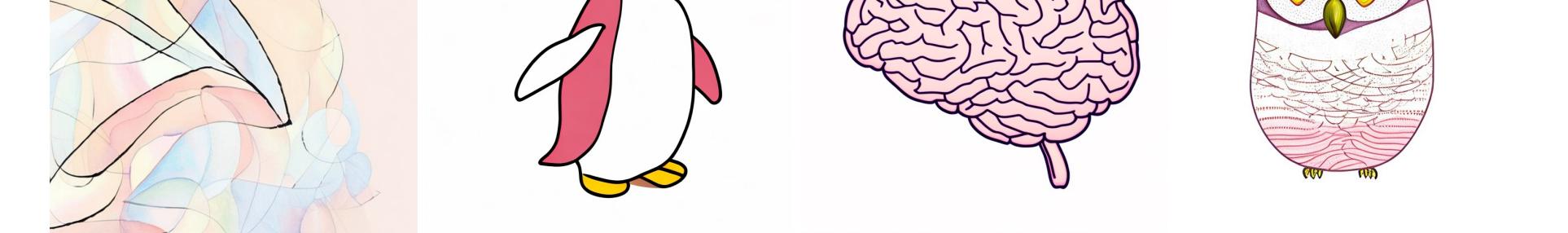


Style-adaptation with the original diffusion model

SD 2.1 [4]



SD 2.1 [4] with ours



SD 2.1 [4]



SD 2.1 [4] with ours

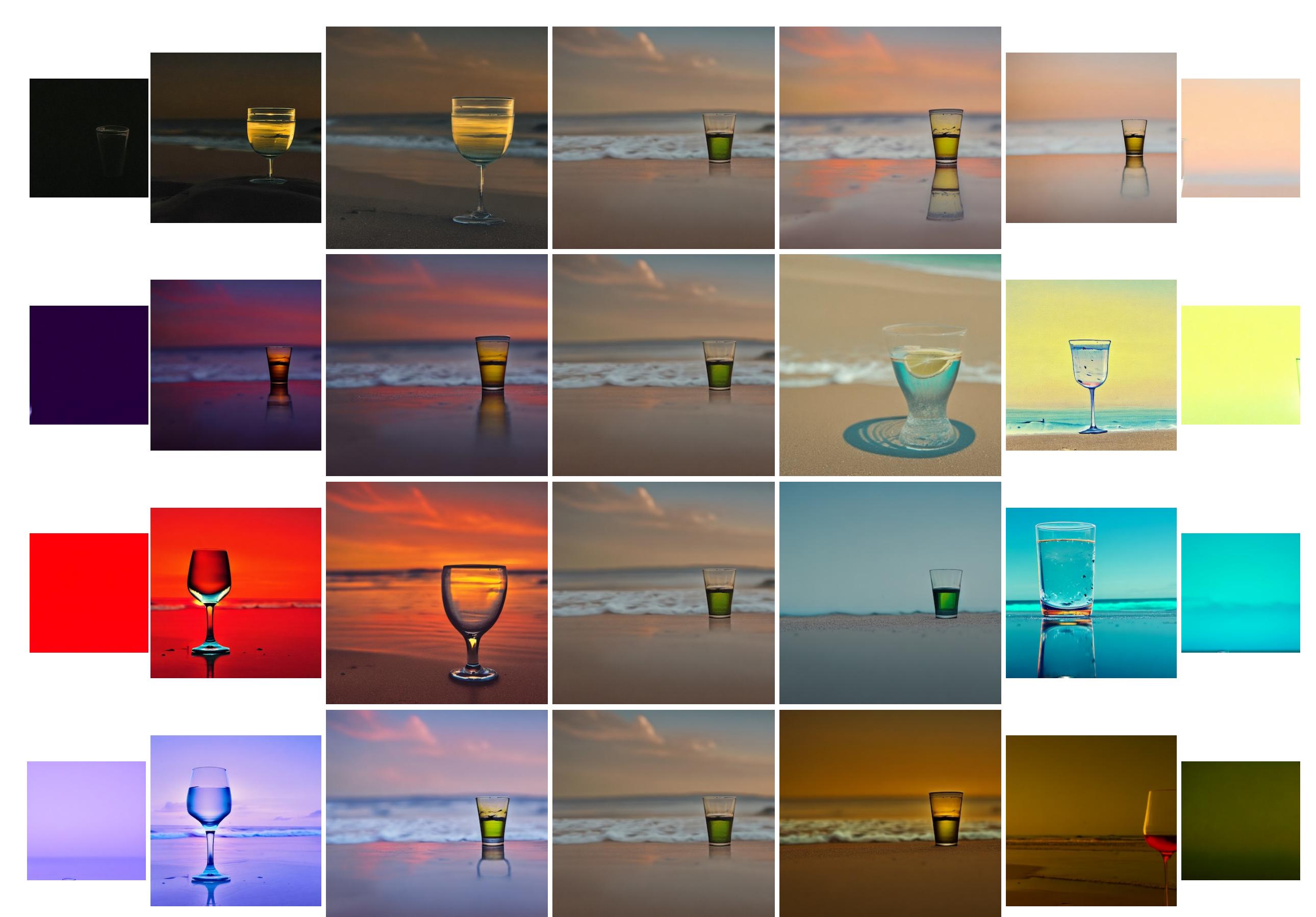


More control on low-frequency components



Setting manually the signal-leak $\sqrt{\bar{\alpha}_T} \tilde{x}$ in \hat{x}_T

→ control on the low-frequency components (e.g., the mean color of the generated images)

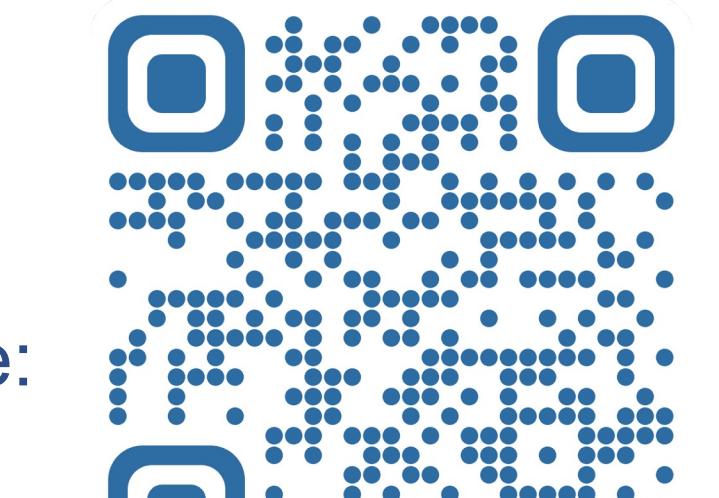


References

"line-art" model [5]: Stable Diffusion v1.4 finetuned with Textual Inversion [5,6] on 7 line-art images [5] (bright background, pastel colors)
"nasa space" model [7]: Stable Diffusion v2 finetuned with DreamBooth [7,8] on 24 photos of astronomical phenomena [7]
Blue city at night: using 9 images from <https://unsplash.com/collections/67793987> (Credits: Unsplash, @borkography)

- [1] Guttenberg, Diffusion with Offset Noise, 2023
- [2] Lin et al., Common Diffusion Noise Schedules and Sample Steps are Flawed, arXiv 2023
- [3] Everaert et al., Diffusion in Style, ICCV 2023
- [4] Stability AI, Stable Diffusion 2.1, 2022 + Rombach et al., High-Resolution Image Synthesis with Latent Diffusion Models, CVPR 2022
- [5] Karan, "line-art" model, <https://huggingface.co/datasets/librariyline-art>, 2022
- [6] Gal et al., An Image is Worth One Word: Personalizing Text-to-Image Generation using Textual Inversion, ICLR 2023
- [7] MatAlart, "nasa space" model, <https://huggingface.co/dreambooth-library/nasa-space-v2>, 2022
- [8] Ruiz et al., DreamBooth: Fine Tuning Text-to-Image Diffusion Models for Subject-Driven Generation, CVPR 2023

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Project website:

<https://ivrl.github.io/signal-leak-bias/>