

CoreML for Stable Diffusion

Analysis and Investigation

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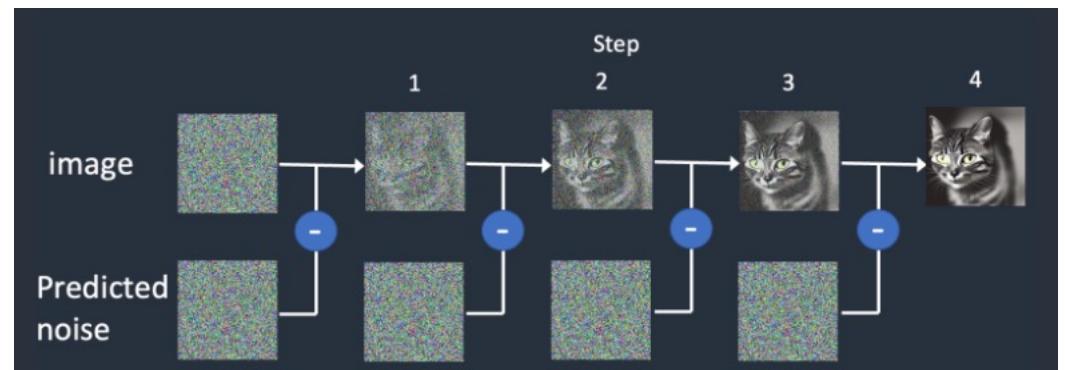
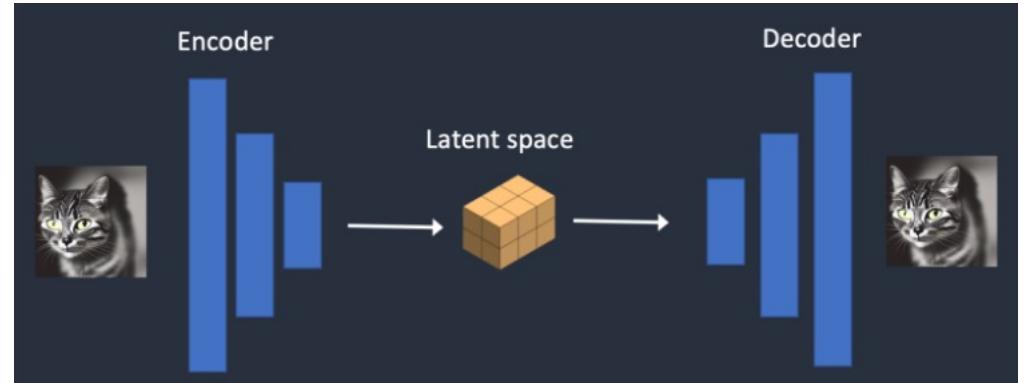
Sept 12, 2023

Overview

- Quick Introduction to Stable Diffusion
- Project Goal and mobile Deployment Pipeline
- Initial tests on macbook
- Deployment to mobile
 - Optimization techniques
 - Testing results
 - SD v1-5
 - SD v2-1
 - SD XL
- Comparisons
- Live Demo
- Conclusion & Next Steps

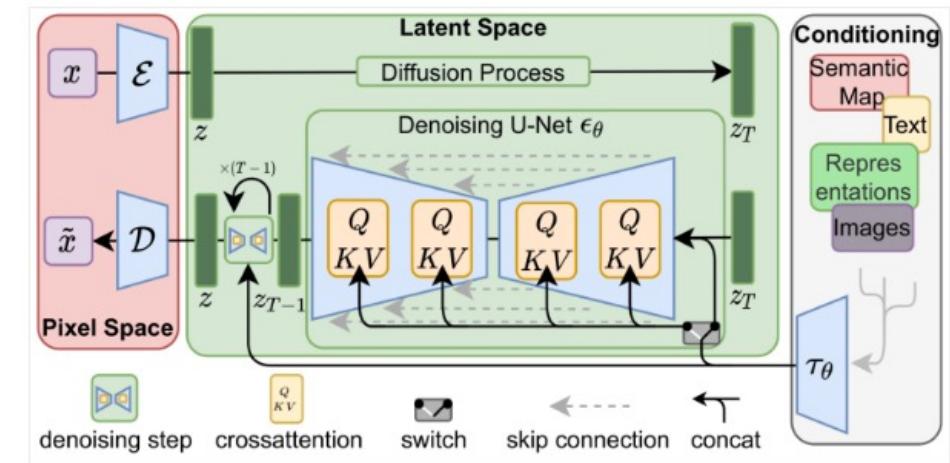
Stable Diffusion – Simplified Explanation

- Latent diffusion model
 - Utilizes a variational autoencoder to compress an image into a smaller latent space
- UNet as noise predictor
 - Generate a random image, add some noise, and have UNet predict the amount of noise
 - For inference, use random noise, and after subtraction will give “generated image”



Stable Diffusion – Simplified Explanation

- Text Conditioning on the UNet
 - Text embeddings fed into UNet via a cross-attention mechanism
 - Network learns to associate latent image features with text embedding features
- Inference:
 - Random noise encoded to latent space
 - Latent noise iteratively subtracted using UNet with text-conditioning
 - Final latent vector decoded to form generated image



The model: Stable Diffusion v1.5

- Text embeddings from OpenAI CLIP ViT-L/14 text-encoder
- Training:
 - 595,000 steps from v1.2 checkpoint
 - LAION-aesthetics v1 5+ dataset, originally on LAION-5B
 - 10% dropping of text-conditioning

Project Goal



A screenshot of a web interface from prompthero.com showing an AI-generated image of a tall ship at sunset. The image is highly detailed and realistic, with warm sunlight streaming through the sails. Below the image is a text box containing a complex AI prompt. A large blue arrow points from this interface to a smartphone screen on the right.

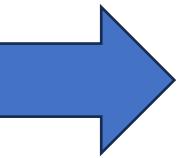
Photo of an ultra realistic sailing ship, dramatic light, pale sunrise, cinematic lighting, low angle, trending on artstation, 4k, hyper-realistic, focused, extreme details, cinematic, masterpiece, intricate artwork by john william turner

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PromptHero Academy
Learn to create breathtaking AI art images

Want to learn how to create images like this one?
Check out our [crash course in prompt engineering & AI art generation!](#)

Fray posted 25 days ago
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My iPhone – specs

- iPhone 13 pro max



A15 Bionic chip

New 6-core CPU with 2 performance and 4 efficiency cores

New 5-core GPU

New 16-core Neural Engine

- 6 GiB of RAM

Deployment pipeline

Following <https://github.com/apple/ml-stable-diffusion>

- Install repository and dependencies
- Download SD model checkpoints (pytorch)
- Convert to Core ML model files (.mlpackage)
- Deploy models on iPhone (iOS 17-beta) – using xCode 15-beta
- Deploy model using apple's StableDiffusion library in Swift, and achieve optimization with CPU + NeuralEngine

Initial Exploration

- Logbook and notes at Notion site: <https://stump-milkshake-736.notion.site/Stable-Diffusion-Mobile-Generation-54bdfc96383f45d7992d164ea62b38ab?pvs=4>
- First tried to run SD model on my MacBook Pro (M1)

Running SDv1-5 on MacBook M1

"An image of a squirrel in Picasso style"



"Macro photography of dewdrops on a spiderweb"



"Underwater photography of a coral reef, with diverse marine life and a scuba diver for scale"



Many ways to run on Mac

- Hugging face diffusers pipeline (python)
- Apple ml-stable-diffusion swift pipeline
- Image generation takes around 0.6s per iteration

"a photo of an astronaut riding
a horse on mars"



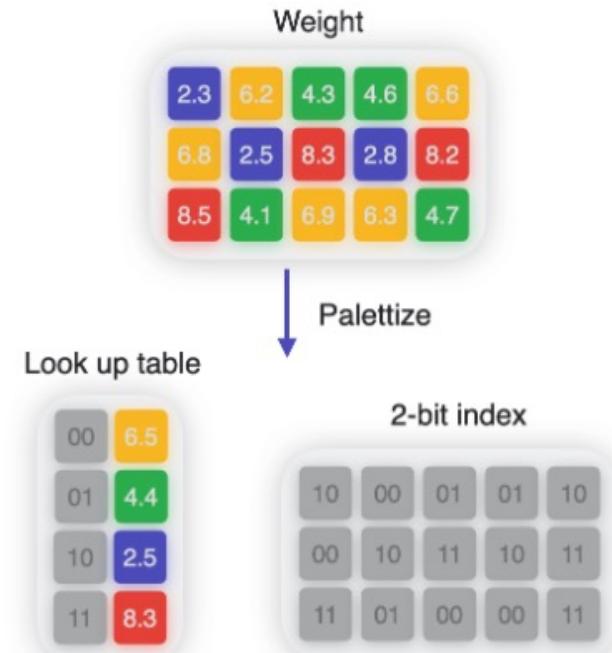
Moving on to mobile

- Apple recommends techniques for optimizing models for deployment on iPhone/iPad
- Very memory intensive (only 6GiB RAM on iPhone 13 pro)
- After initial exploration:
 - Must update to iOS 17 beta on iPhone
 - Built custom app using Xcode 15 beta



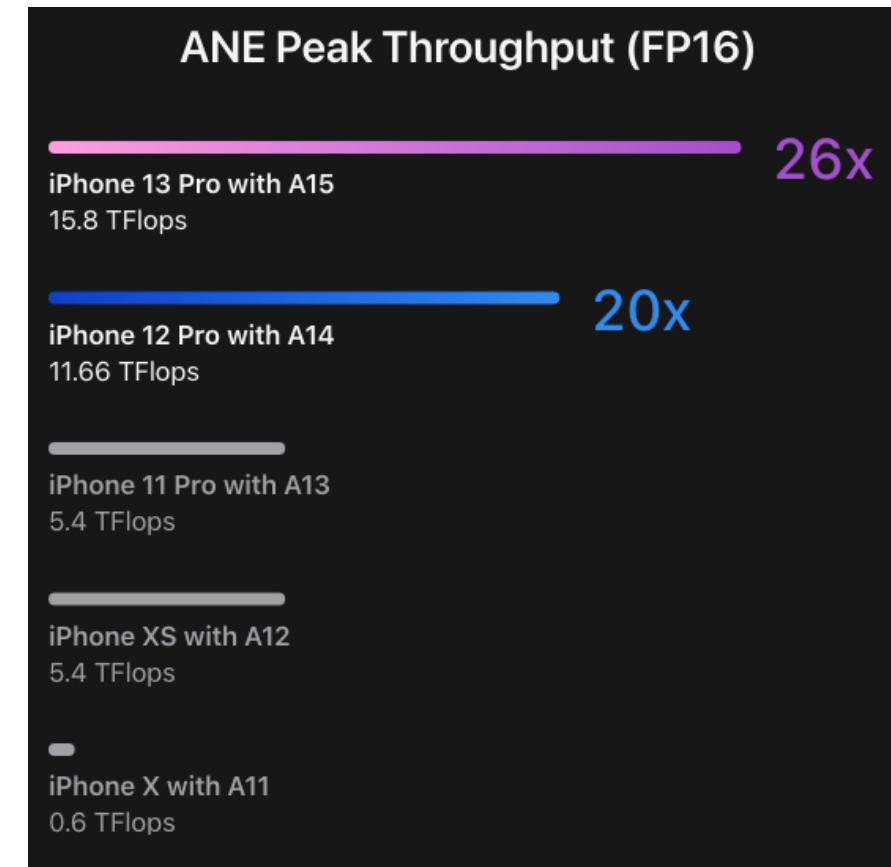
Palettization technique

- Clusters weights in model to a lookup table
- Reduces size of weights.
- Decompressing palettized weights happen “just in time” on iOS 17 +, leading to enhanced latency



Accelerating Transformers with NeuralEngine

- Apple Neural Engine (ANE)
 - Specialized operations on Tensors to enhance performance
- Chunks input tensors
- Use batched matrix multiplication (einsum formula) to avoid extra memory copying



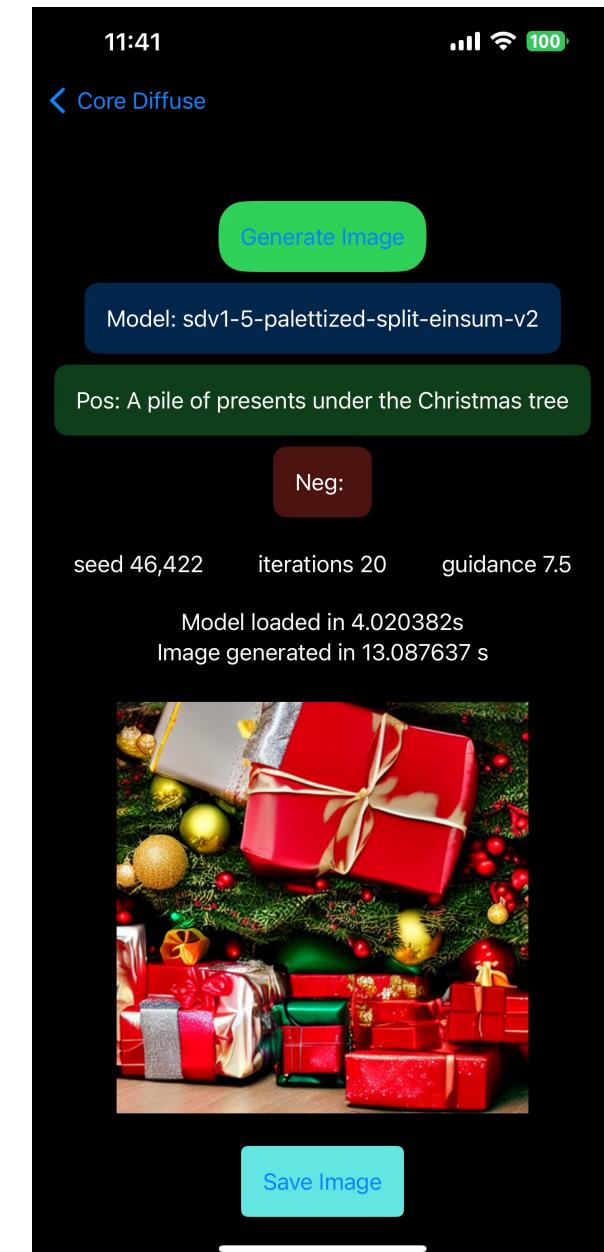
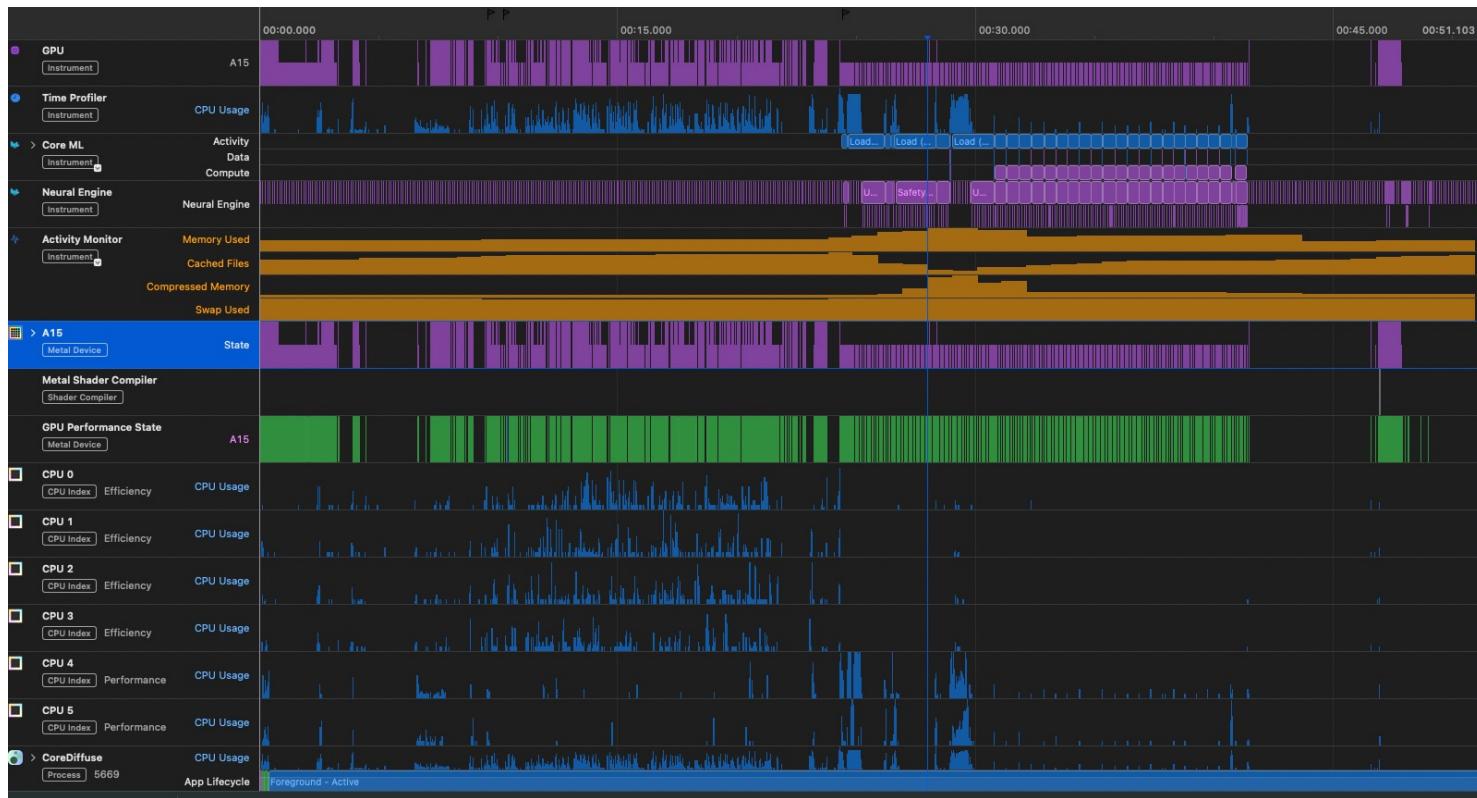
Stable diffusion 1.5 model

- Total size 0.957 GiB
- 6 bit palettization
- Using split-einsum v2

	sdv1-5-palettized-split-einsum-v2	--
	merges.txt	525 KB
	TextEncoder.mlmodelc	140.1 MB
	Unet.mlmodelc	648.2 MB
	VAEDecoder.mlmodelc	99.2 MB
	VAEEncoder.mlmodelc	68.5 MB
	vocab.json	862 KB

Generation of one image using sdv1-5

- Peak memory usage: 5.04 GiB
- Peak CPU usage: 440%

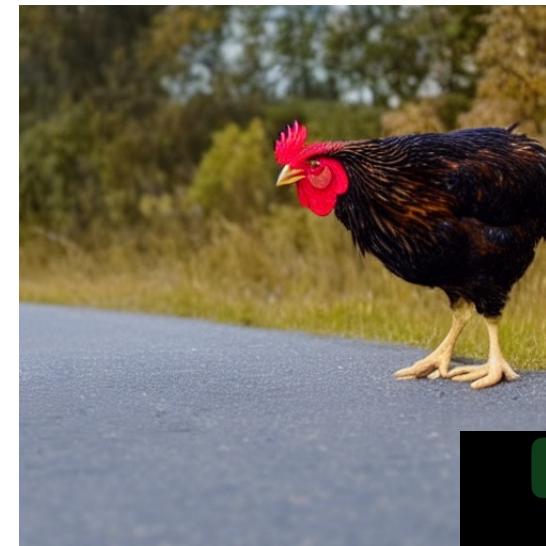


Sdv1-5 image generation

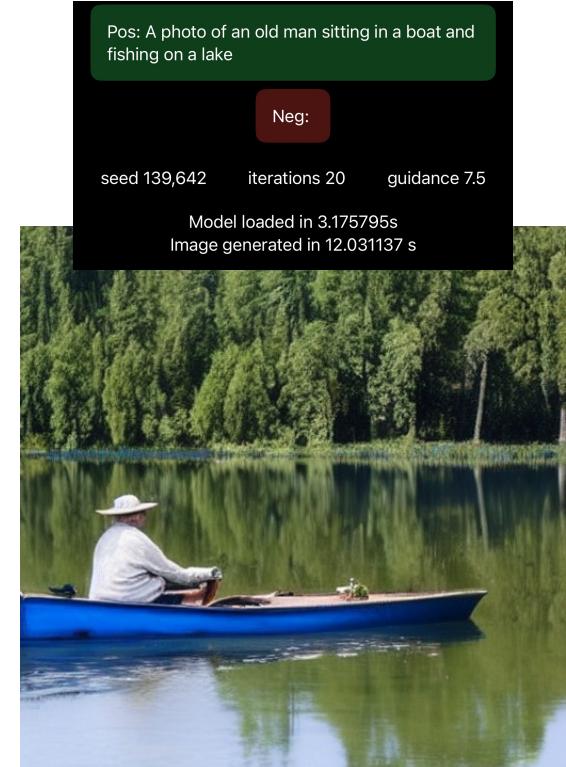
- Initial loading of model takes around 120s
- Afterwards, model loading takes around 3.5s
- Image generation takes ~ 0.75s / step



Pos: A cartoon drawing of a dragon guarding a pile of treasure
Neg:
seed 139,642 iterations 20 guidance 7.5
Model loaded in 3.298039s
Image generated in 14.491954 s



Pos: A chicken standing on the road
Neg:
seed 139,642 iterations 20 guidance 7.5
Model loaded in 3.239590s
Image generated in 12.299682 s



Pos: A photo of an old man sitting in a boat and fishing on a lake

Neg:

seed 139,642 iterations 20 guidance 7.5

Model loaded in 3.175795s
Image generated in 12.031137 s

Stable Diffusion 2.1 model

- Total size 1.14 GiB
- 6 bit palettization
- Using split-einsum v2

▼	sdv2-1-palettized-split-einsum-v2	--
	merges.txt	525 KB
	TextEncoder.mlmodelc	319.3 MB
	Unet.mlmodelc	653.1 MB
	VAEDecoder.mlmodelc	99.2 MB
	VAEEncoder.mlmodelc	68.5 MB
	vocab.json	862 KB

V2.1 stats

- Peak memory usage 5.06 GiB
- Peak CPU usage around 420%

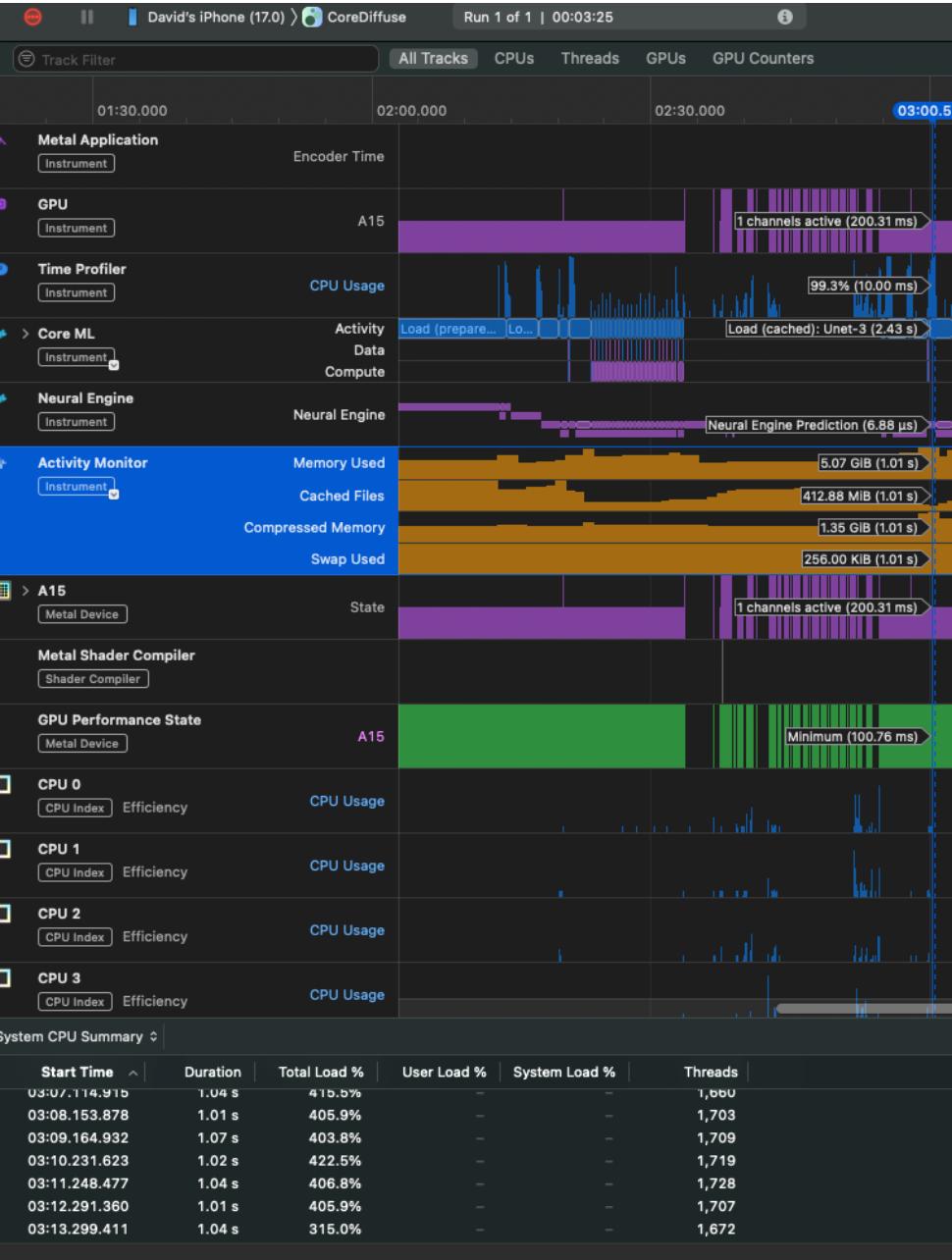
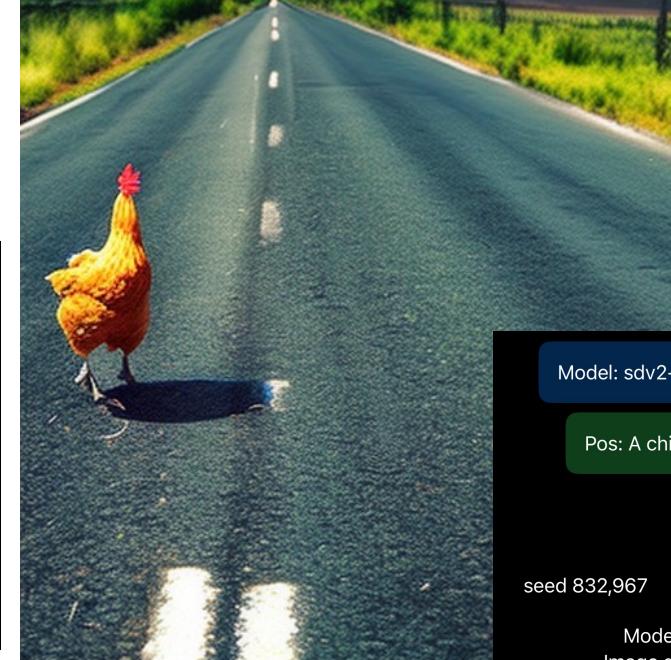


Image Generation with v2-1

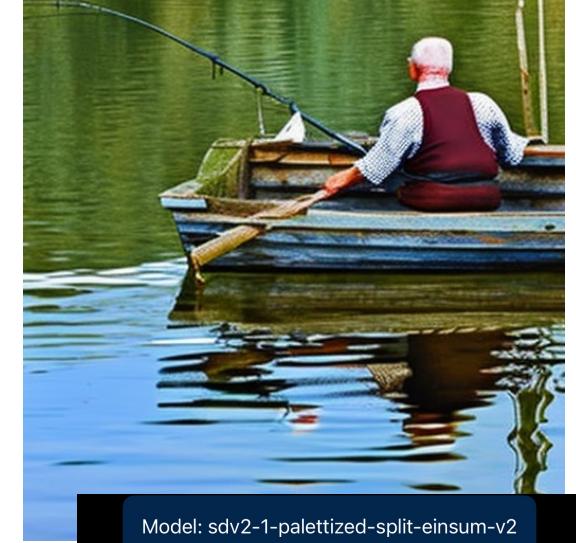
- Model loading around 120s for the first time
- Successive loading takes around between 1.5 to 4.0s
- Image generation takes around 0.75s / step (but slightly faster than 1.5)



Model: sdv2-1-palettized-split-einsum-v2
Pos: A cartoon drawing of a dragon guarding a pile of treasure
Neg:
seed 139,642 iterations 20 guidance 7.5
Model loaded in 2.008189s
Image generated in 12.390302 s



Model: sdv2-1-palettized-split-einsum-v2
Pos: A chicken standing on the road
Neg:
seed 832,967 iterations 20 guidance 7.5
Model loaded in 3.058700s
Image generated in 11.258951 s



Model: sdv2-1-palettized-split-einsum-v2
Pos: A photo of an old man sitting on a boat and fishing on a lake
Neg:
seed 832,967 iterations 20 guidance 7.5
Model loaded in 2.277595s
Image generated in 12.667799 s

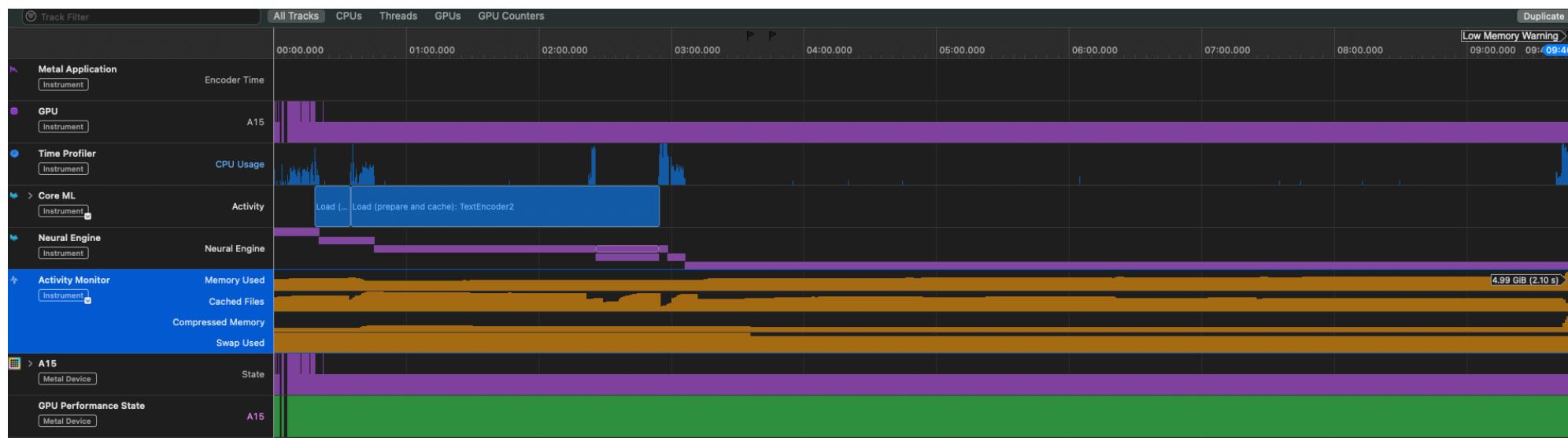
SDXL models

- Total model size 3.36 GiB
- Not yet supported in CoreML
- Options for 6bit, 4.5bit, and 3.6bit palettization

sdxl-base-mbp-4-50-palettized	
merges.txt	525 KB
TextEncoder.mlmodelc	246.3 MB
TextEncoder2.mlmodelc	1.39 GB
Unet.mlmodelc	1.45 GB
VAEDecoder.mlmodelc	198.1 MB
VAEEncoder.mlmodelc	68.5 MB
vocab.json	862 KB

```
E5RT encountered an STL exception. msg = MILCompilerForANE error: failed to compile ANE model  
using ANEF. Error=_ANECompiler : ANECCCompile() FAILED.
```

```
E5RT: MILCompilerForANE error: failed to compile ANE model using ANEF. Error=_ANECompiler :  
ANECCCompile() FAILED (11)
```

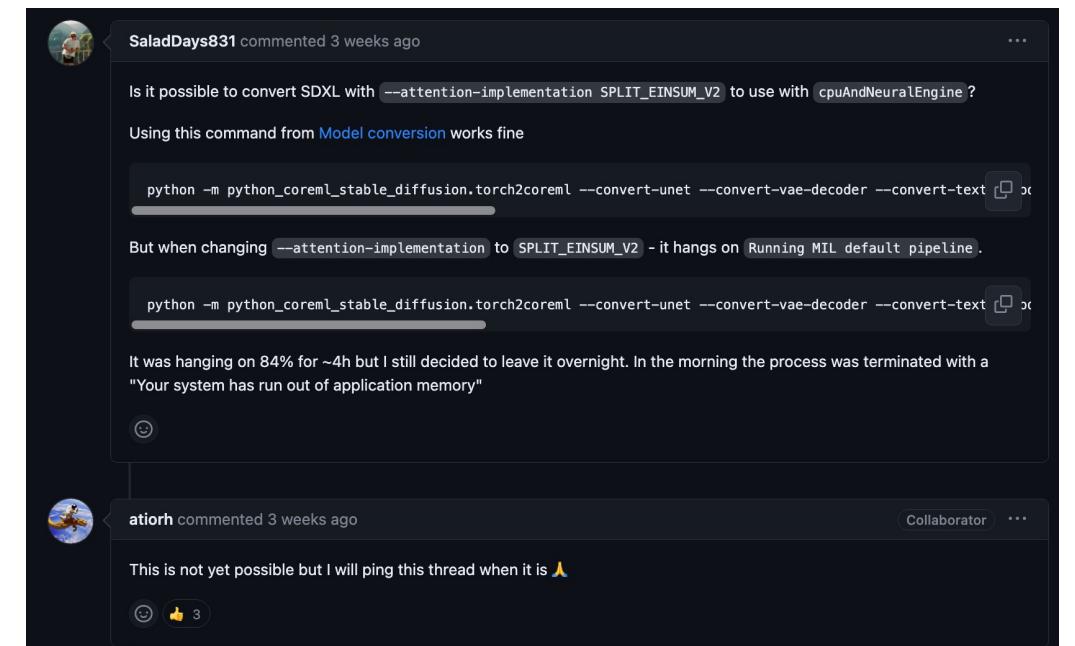
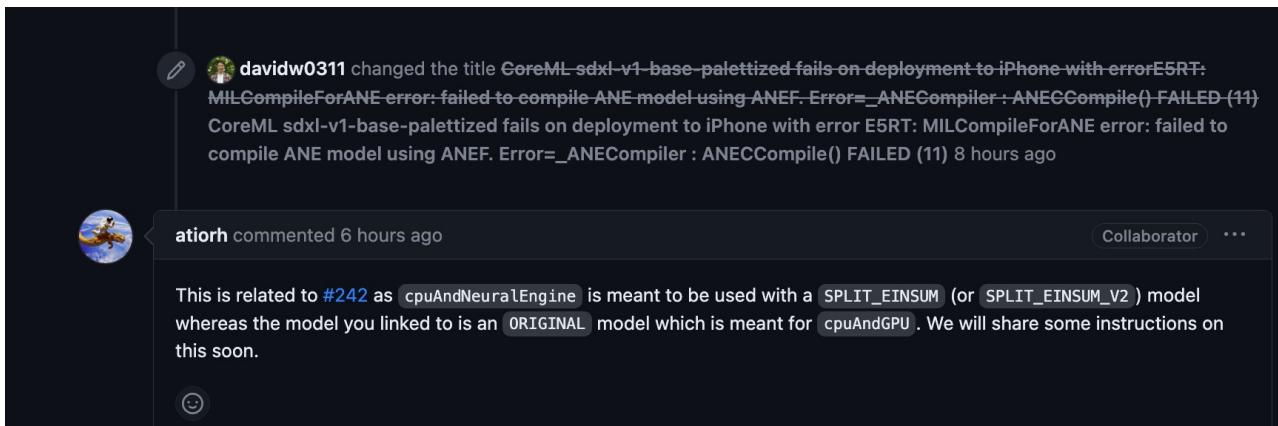


Model crashes

- <https://github.com/apple/ml-stable-diffusion/issues/228>
 - Issue not yet resolved, the apple coreml team is currently working on the official release
 - <https://github.com/apple/ml-stable-diffusion/issues/255>
- To achieve optimization using CoreML (on Mac) requires upgrading to OS 14 beta

Raised error to ml-stable-diffusion github

- <https://github.com/apple/ml-stable-diffusion/issues/255>
 - Response: split-einsum conversion for xl models is not currently supported
 - Apple team is currently working on resolving the issue, should be available soon



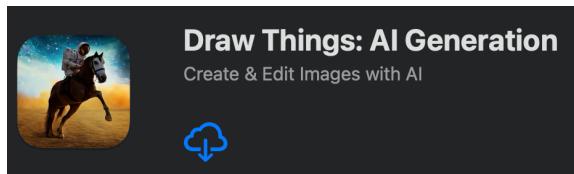
V1.5 vs V2.1

- Similar speed of image generation, v2.1 slightly faster
- V2.1 model size larger (1.14 GiB vs 0.957 GiB)
- V2.1 performs better with negative prompts
- Both only support 512x512 px image generation (for now)



1.5 and 2.1 model performance still very fast

- Comparison to Draw Things App

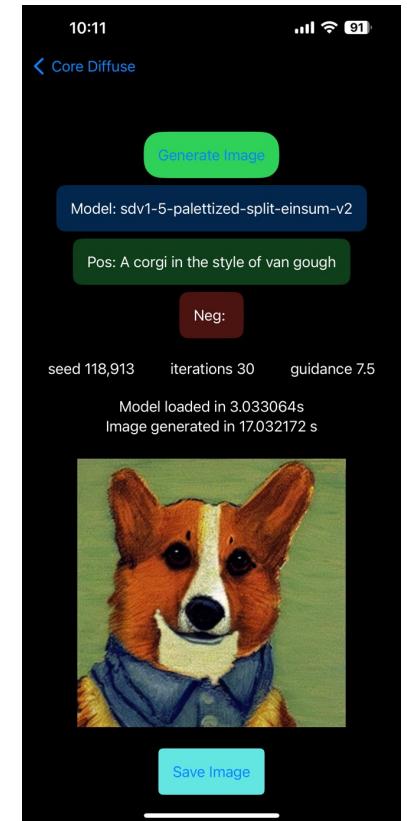


30 steps of image generation on sd1.5 takes ~ 80s



30 steps of image generation on sd1.5 takes ~ 20s

4x speedup!



Reflection

- Challenging project
 - Lack of resources/documentation on newest CoreML features
 - No prior experience with Swift or Apple app development
 - Very early-stage development, only compatible with iOS 17 beta, xcode 15 beta, and OS 14 beta
 - Lack of storage space on Mac after trying and downloading many models
- Deployed and investigated performance of stable diffusion v1-5, v2-1, and xl models on iPhone 13 pro max, accelerated with coreml and apple neural engine

Next Steps / Future work

- Investigate crash errors of loading xl model on coreml, and attempt to resolve
- Investigate full memory usage of running models, and further optimize performance.
- Extensively compare performance between different palettized models.
- Integrate LoRA checkpoints onto of sd models, optimized via coreml
- Better app UI and deployment
- Allow use of control-net

Goal?



< Core Diffuse

Generate Image

Model: sdv2-1-palettized-split-einsum-v2

Pos: Photo of an ultra realistic sailing ship, dramatic light, pale sunrise, cinematic lighting, low angle, trending on artstation, 4k, hyper-realistic, focused, extreme details, cinematic, masterpiece, intricate artwork by john william turner

Neg: Blurry, unclear, low resolution

seed 41,847 iterations 50 guidance 7.5

Model loaded in 3.649035s

Image generated in 24.040052 s



Project repo:

<https://github.com/davidw0311/CoreDiffuse>

References

- <https://github.com/apple/ml-stable-diffusion>
- <https://github.com/huggingface/swift-coreml-diffusers>
- <https://github.com/madebyollin/maple-diffusion>
- <https://github.com/ynagatomo/ImgGenSD2>
- <https://jalammar.github.io/illustrated-stable-diffusion/>
- <https://liuliu.me/eyes/stretch-iphone-to-its-limit-a-2gib-model-that-can-draw-everything-in-your-pocket/>
- <https://arxiv.org/pdf/2112.10752.pdf>
- <https://machinelearning.apple.com/research/stable-diffusion-coreml-apple-silicon>