# 50.039 Theory and Practice of Deep Learning W12-S2 About Diffusion Models

Matthieu De Mari



# About this week (Week 12, an informative lecture about diffusion models)

- 1. What is a **Stochastic Differential Equation**?
- 2. What are typical uses of SDEs?
- 3. Can I noise images progressively using SDEs?
- 4. Can you train a denoising autoencoder model that removes SDE noise from images?
- 5. What is the **diffusion process** and how to use it to train **diffusion models**?
- 6. What are the architectures of DALL-E and Midjourney models?

#### Diffusion Models: a definition

#### **Definition (Diffusion Models):**

Diffusion models are a class of stochastic mathematical models that have typically been used to describe the random motion of particles, often referred to as "Brownian motion."

They are widely used in various fields, such as physics, chemistry, biology, and finance, for instance to model the stochastic behaviour of particles or other quantities that exhibit random fluctuations over time.

The fundamental concept behind diffusion models is the idea of **Stochastic Differential Equations (SDEs)**.

## Stochastic Differential Equations (SDEs)

#### **Definition (Stochastic Differential Equations):**

Stochastic Differential Equations (SDEs) are differential equations where one or more of the terms include a random component, usually in the form of a noise term.

The randomness is often modelled using a Wiener process or Brownian motion, which is a continuous-time stochastic process characterized by a mean, a variance, and independent increments.

#### Diffusion models as SDEs

Mathematically, a diffusion model can be represented as an SDE:

$$dx(t) = a(x(t),t)dt + b(x(t),t)dW(t)$$
 and  $x(t=0) = x_0$ 

Where,

- x(t) represents the state of the system at time t,
- a(x(t), t) is the **deterministic drift term**,
- b(x(t), t) is the stochastic diffusion term,
- and W(t) is the Wiener process.

The drift term governs the deterministic part of the system's dynamics, while the diffusion term introduces randomness into the system's evolution.

#### A first SDE: the Brownian motion

As a first example, let us consider one of the simplest SDEs:

$$dx(t) = dW(t)$$
 and  $x(0) = 0$ 

This SDE typically defines a random walk starting from 0 and the noise process is simply zero mean and variance  $\sqrt{timestep}$ .

Implemented as shown in Notebook 1 and on the right.

```
1 # Parameters
2 num_simulations = 5
3 total_time = 10.0
4 num_steps = 1000
5 dt = total_time/num_steps
```

```
# Initialize the state variable as 0, at t = 0
x = np.zeros((num_simulations, num_steps))
x[:, 0] = 0.0

# Simulate Brownian motion
for t in range(1, num_steps):
    # Definition random motion using a Normal distribution with
# Zero mean and variance sqrt(dt)
dW = np.random.normal(0, np.sqrt(dt), size = num_simulations)
x[:, t] = x[:, t-1] + dW
```

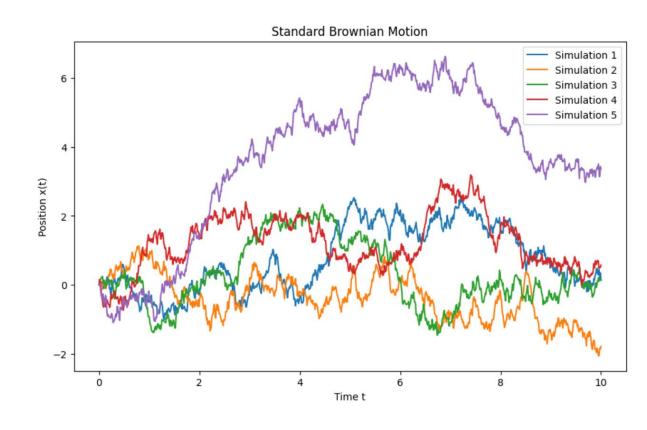
#### A first SDE: the Brownian motion

As a first example, let us consider one of the simplest SDEs:

$$dx(t) = dW(t)$$
 and  $x(0) = 0$ 

This SDE typically defines a random walk starting from 0 and the noise process is simply zero mean and variance  $\sqrt{timestep}$ .

Implemented as shown in Notebook 1 and on the right.



#### A second SDE with drift

As a second example, let us consider an SDE wit drift:

$$dx(t) = 0.1x(t)dt + dW(t)$$
 and  $x(0) = 10$ 

The drift function here is 0.1x(t), adding deterministically ~10% to the current value of x(t) for each unit of time t.

#### A second SDE with drift

- This second SDE is typically implemented as shown.
- We also show as xr, the realisation of the SDE without any Wiener motion.

```
1 # Parameters
2 num_simulations = 5
3 total_time = 10.0
4 num_steps = 1000
5 dt = total_time / num_steps
```

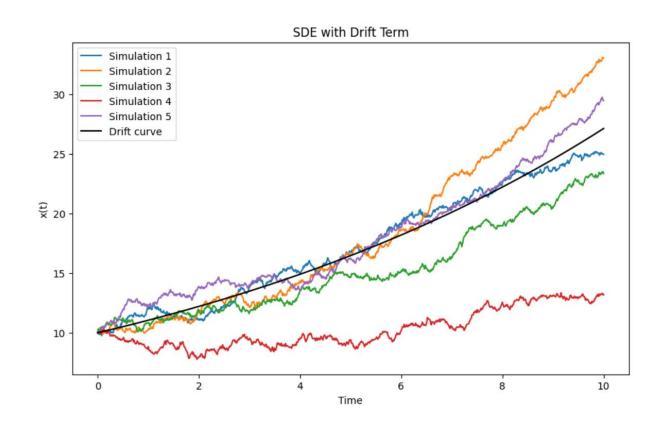
```
1 # Define the drift term
2 def drift(x, t):
3 return 0.1*x
```

```
# Initialize the state variable
x = np.zeros((num_simulations, num_steps))
x[:, 0] = 10.0 # Initial condition

# Simulate the SDE
for t in range(1, num_steps):
    dW = np.random.normal(0, np.sqrt(dt), size = num_simulations)
    dx = drift(x[:, t-1], t*dt)*dt + dW
    x[:, t] = x[:, t-1] + dx
```

#### A second SDE with drift

- This second SDE is typically implemented as shown.
- We also show as xr, the realisation of the SDE without any Wiener motion.
- Typically a 10% increase per unit of time (as is often described the US stock market), with noisy variations.
- Is the stock market an SDE of some sort then?



## Typical examples of SDEs

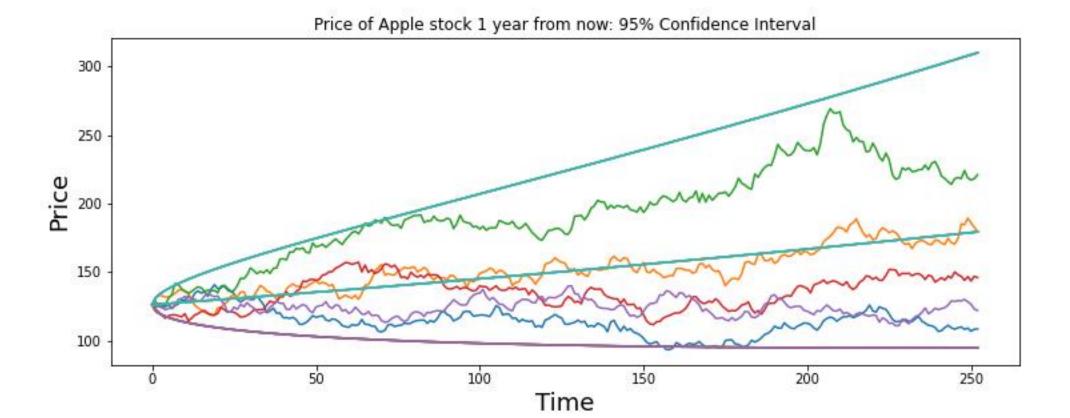
In fact, stochastic differential equations (SDEs) have been widely used, in finance, to model the dynamics of various financial instruments and quantities. For instance, the Geometric Brownian motion is a popular model for simulating the price of assets such as stocks. It is given by:

$$dS(t) = \mu S(t)dt + \sigma S(t)dW(t)$$

Where S(t) is the asset price,  $\mu$  is the drift of the asset (i.e. expected return of the asset over time),  $\sigma$  is the volatility, and W(t) is a Wiener process modelling the unpredictable behaviour of the asset.

# Typical examples of SDEs

In fact, stochastic differential equations (SDEs) have been widely used, in finance, to model the dynamics of various financial instruments and



- **Task:** take images from MNIST and noise them. The model should try its best to denoise them an get back to the original image.
- Dataset: MNIST images, noised (inputs) and original (output).
- **Model:** an encoder-decoder model of some sort, similar to an autoencoder (almost).
- Loss and training procedure: As in the case of autoencoders, a simple MSE loss and a basic training procedure with Adam optimizer will do.

 For simplicity, we consider an encoder-decoder model, similar to our autoencoder in W10-NB1, running on linear layers.

```
class Denoiser(nn.Module):
    def __init__(self):
        super(Denoiser, self).__init__()
        self.encoder = Encoder()
        self.decoder = Decoder()

def forward(self, x):
    return self.decoder(self.encoder(x))
```

```
# Define the encoder
   class Encoder(nn.Module):
       def init (self):
           super(Encoder, self).__init__()
            self.encoder = nn.Sequential(
                nn.Linear(28*28, 400),
               nn.ReLU(),
                nn.Linear(400, 400),
 8
               nn.ReLU()
10
11
       def forward(self, x):
12
13
            return self.encoder(x)
```

```
# Define the decoder
   class Decoder(nn.Module):
       def init_(self):
           super(Decoder, self). init ()
           self.decoder = nn.Sequential(
               nn.Linear(400, 400),
               nn.ReLU(),
               nn.Linear(400, 28*28),
               nn.Sigmoid()
 9
10
11
       def forward(self, x):
12
           return self.decoder(x)
13
```

- For simplicity, we consider an encoder-decoder model, similar to our autoencoder in W10-NB1, running on linear layers.
- And a basic training procedure...

```
class Denoiser(nn.Module):
    def __init__(self):
        super(Denoiser, self).__init__()
        self.encoder = Encoder()
        self.decoder = Decoder()

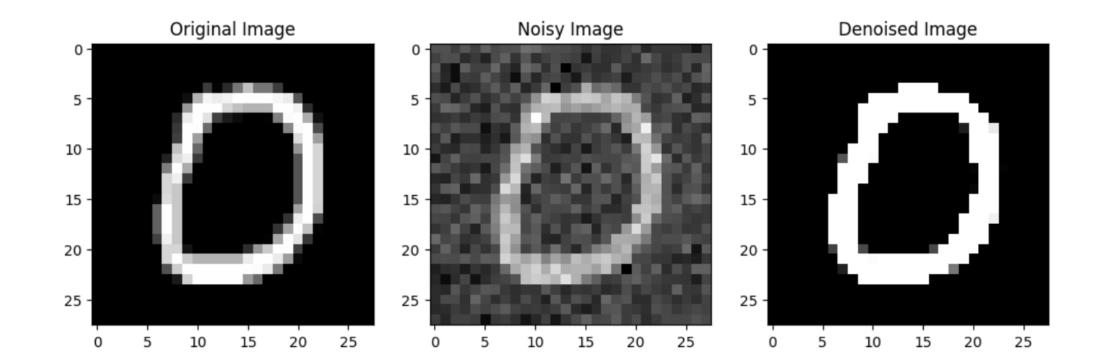
def forward(self, x):
    return self.decoder(self.encoder(x))
```

```
# Define the encoder
   class Encoder(nn.Module):
       def init (self):
           super(Encoder, self).__init__()
           self.encoder = nn.Sequential(
                nn.Linear(28*28, 400),
               nn.ReLU(),
                nn.Linear(400, 400),
 8
               nn.ReLU()
10
11
12
       def forward(self, x):
13
            return self.encoder(x)
```

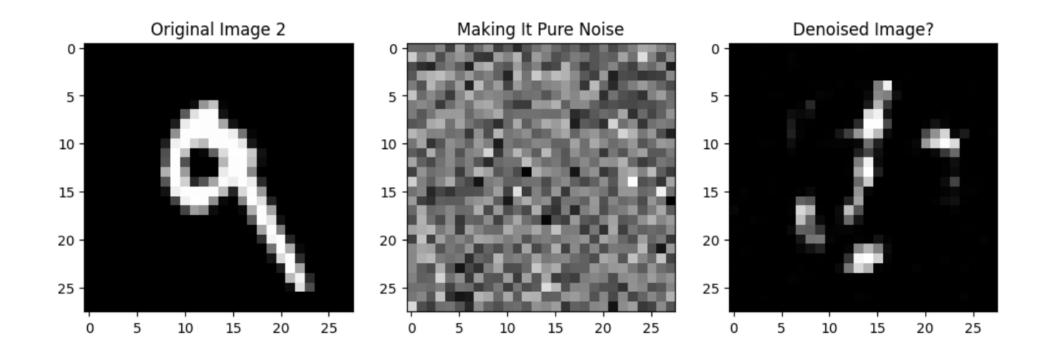
```
# Define the decoder
   class Decoder(nn.Module):
       def init_(self):
            super(Decoder, self).__init__()
            self.decoder = nn.Sequential(
                nn.Linear(400, 400),
               nn.ReLU(),
                nn.Linear(400, 28*28),
                nn.Sigmoid()
 9
10
11
       def forward(self, x):
12
            return self.decoder(x)
13
```

```
# Hyperparameters
2 epochs = 10
 3 learning rate = 1e-3
4 noise factor = 0.5
1 # Initialize the model, loss function, and optimizer
 2 model = Denoiser()
 3 criterion = nn.MSELoss()
   optimizer = optim.Adam(model.parameters(), lr = learning rate)
   # Training Loop
   for epoch in range(epochs):
       for batch_idx, (data, _) in enumerate(train_loader):
           # Noise the data with one or many steps of SDEs
           # (Here a single step Brownian Motion)
           data = data.view(data.size(0), -1)
 6
           noisy data = data + noise factor*torch.randn like(data)
 8
           # Forward pass
 9
           optimizer.zero grad()
10
           outputs = model(noisy data)
11
12
13
           # Backprop
14
           loss = criterion(outputs, data)
           loss.backward()
15
           optimizer.step()
16
17
18
           # Display
19
           if batch idx % 100 == 0:
20
               print(f'Epoch: {epoch+1}/{epochs}, Batch: {batch_idx}/{len(train_loader)}, Loss: {loss.item()}')
```

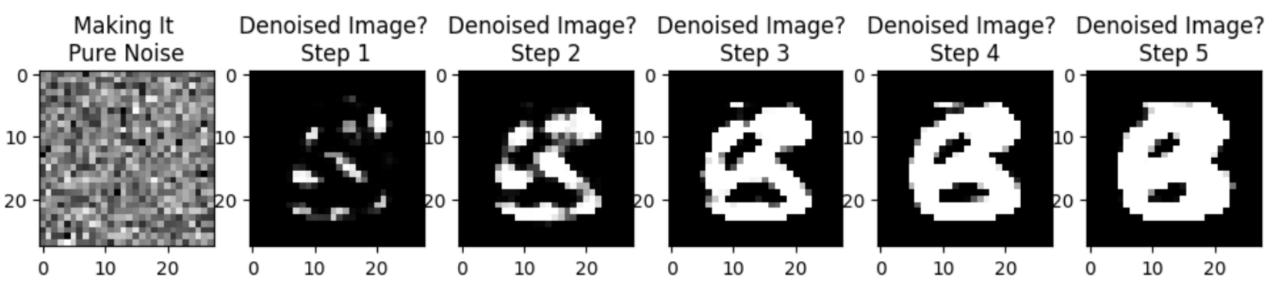
• The model trains and seems able to denoise lightly noised images.



- The model trains and seems able to denoise lightly noised images.
- Could it denoise a pure noise image, though? Not really.

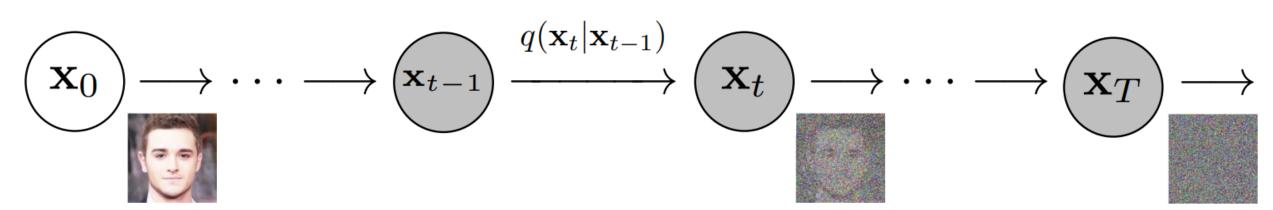


- The model trains and seems able to denoise lightly noised images.
- Could it denoise a pure noise image, though? Not really.
- What if we used the model many (e.g. 5) times in a row on our full noise image? Eventually would obtain something B&W in the likes of MNIST but somehow non-sensical in terms of shapes...?!



**Idea:** what if instead of noising the image in one step, and denoising it in one call of the model, we try iterative noising and denoising?

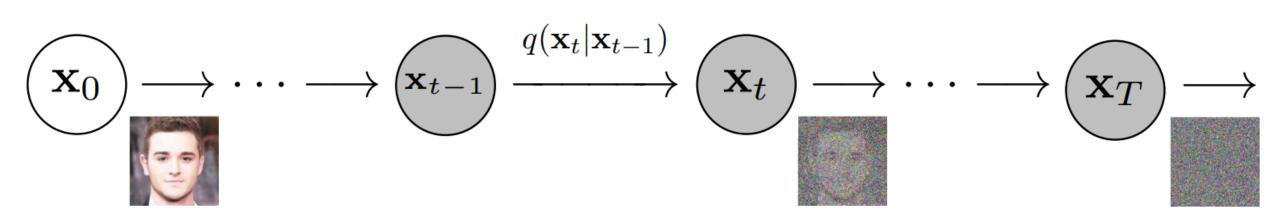
(Somehow similar to the idea of modifying the image iteratively during the iterated versions of our attacks in W5!)



......

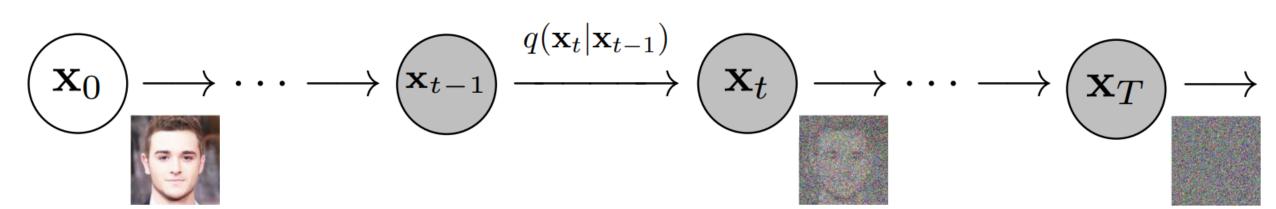
**Idea:** what if instead of noising the image in one step, and denoising it in one call of the model, we try iterative noising and denoising?

- **Step 1:** Add progressive noising T times in a row, until the image becomes pure noise (or so).
- The noising process follows some sort of an SDE equation describing the progressive noising of the original image. (our Forward here!)



**Idea:** what if instead of noising the image in one step, and denoising it in one call of the model, we try iterative noising and denoising?

- Step 2: Train a single model to try and cancel each round of noising.
- It receives image  $x_t$  and tries to produce  $x_{t-1}$ , do for each t. (our Backward here!)



**Idea:** what if instead of noising the image in one step, and denoising it in one call of the model, we try iterative noising and denoising?

- Training procedure: Use original images from MNIST and noise them T times in a row, train model on denoising each step and pair of images  $(x_t, x_{t-1})$ .
- Use MSE loss and training procedure similar as before.

### A basic Diffusion Model setup

- Roughly same encoder-decoder model as before.
- Adds an info about current time step t for noising as input parameter, along with  $x_t$ .

```
class DiffusionModel(nn.Module):
        def __init__(self):
            super(DiffusionModel, self). init ()
            self.network = nn.Sequential(
                nn.Linear(28 * 28 + 1, 400),
                nn.ReLU(),
                nn.Linear(400, 400),
                nn.ReLU(),
                nn.Linear(400, 28 * 28)
10
11
12
        def forward(self, x, t):
13
            x = torch.cat([x, t.expand(x.size(0), 1)], dim = 1)
            return self.network(x)
14
```

### A basic Diffusion Model setup

- Roughly same encoder-decoder model as before.
- Adds an info about current time step t for noising as input parameter, along with  $x_t$ .
- Same training parameters as before.
- Five steps of noising (T = 5).
- Training procedure somewhat similar to what we had before...

```
class DiffusionModel(nn.Module):
        def __init__(self):
            super(DiffusionModel, self). init ()
            self.network = nn.Sequential(
                nn.Linear(28 * 28 + 1, 400),
                nn.ReLU(),
                nn.Linear(400, 400),
                nn.ReLU(),
                nn.Linear(400, 28 * 28)
10
11
12
       def forward(self, x, t):
13
            x = torch.cat([x, t.expand(x.size(0), 1)], dim = 1)
            return self.network(x)
14
```

```
# Hyperparameters
epochs = 10
batch_size = 64
learning_rate = 1e-3
num_timesteps = 5
noise_schedule = torch.linspace(0, 1, num_timesteps)
```

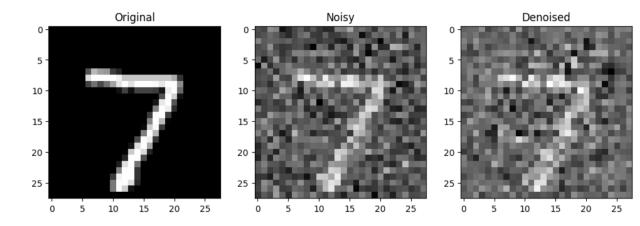
```
# Initialize the model, loss function, and optimizer
model = DiffusionModel()
criterion = nn.MSELoss()
optimizer = optim.Adam(model.parameters(), lr = learning_rate)
```

```
1 # Training Loop
   # (Very long, use at your own discretion!)
   for epoch in range(epochs):
       for batch idx, (data, ) in enumerate(train loader):
           data = data.view(data.size(0), -1)
           noisy data = data.clone()
           # Apply multiple noising steps
           for t in range(num timesteps-1):
                noisy data += torch.randn like(noisy data)*torch.sqrt(noise schedule[t+1]-noise schedule[t])
10
11
12
           optimizer.zero grad()
13
           loss = 0
14
15
           # Apply multiple denoising steps
           for t in reversed(range(num_timesteps-1)):
16
                denoised_data = noisy_data-torch.randn_like(noisy_data)*torch.sqrt(noise_schedule[t+1]-noise_schedule[t])
17
               target_score = (data-denoised_data)/(noise_schedule[t+1]-noise_schedule[t])
18
                score_prediction = model(denoised_data, noise_schedule[t].view(-1, 1))
19
20
               loss += criterion(score prediction, target score)
21
           loss.backward()
22
           optimizer.step()
23
24
25
           if batch idx % 100 == 0:
                print(f'Epoch: {epoch+1}/{epochs}, Batch: {batch idx}/{len(train loader)}, Loss: {loss.item()}')
26
```

## In practice, however...

Need more iterations and many more steps of noising and denoising (heavy to train), in the hundreds/thousands or so.

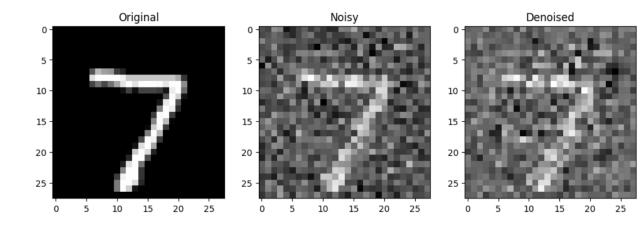
- Our Notebook 3 runs on only 5 steps of noising/denoising.
- But you get the idea...

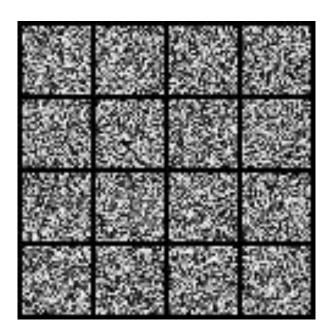


## In practice, however...

# With more steps and CNN layers, get this kind of behaviour:

 Being able to train a model to denoise and eventually obtain a convincing MNIST picture from something that was originally pure noise! (as shown in GIF).

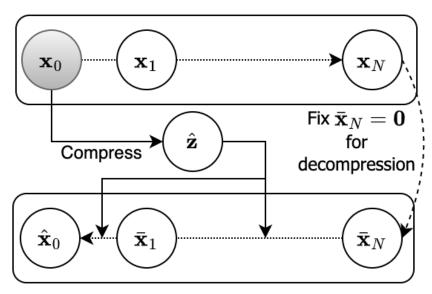




## Making it conditional

- If you train with a dataset of captioned images, the extra label could be used during the denoising part, to help the denoising model figure out what was the original image about.
- Encode these captions and use them during training of the denoising model in place of t.
- Helps to direct the denoising in right direction!

#### Diffusion



#### Conditional Denoising



The man at bat readies to swing at the pitch while the umpire looks on.



A large bus sitting next to a very tall



A horse carrying a large load of hay and two people sitting on it.

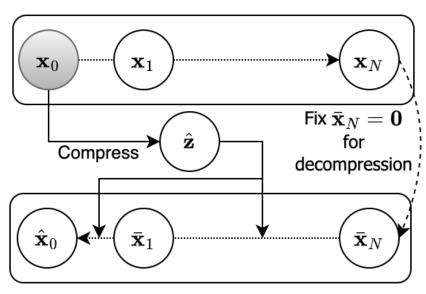


Bunk bed with a narrow shelf sitting underneath it.

# Making it conditional

- During testing, use pure noise as  $x_N$  and ask the user for a caption of their choice that will be compressed as  $\hat{z}$ .
- Use your denoiser many times in a row and obtain an image matching the caption  $\hat{z}$  starting from pure noise?!

#### Diffusion



#### **Conditional Denoising**



The man at bat readies to swing at the pitch while the umpire looks on.



A large bus sitting next to a very tall



A horse carrying a large load of hay and two people sitting on it.

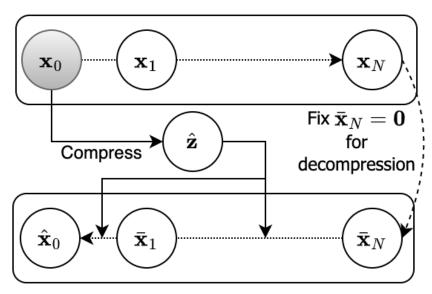


Bunk bed with a narrow shelf sitting underneath it.

# Making it conditional

- During testing, use pure noise as  $x_N$  and ask the user for a caption of their choice that will be compressed as  $\hat{z}$ .
- Use your denoiser many times in a row and obtain an image matching the caption  $\hat{z}$  starting from pure noise?!
- Yes, and better than most GANs these days!

#### Diffusion



#### Conditional Denoising



The man at bat readies to swing at the pitch while the umpire looks on.



A large bus sitting next to a very tall



A horse carrying a large load of hay and two people sitting on it.



Bunk bed with a narrow shelf sitting underneath it.

#### Restricted

#### Stable Diffusion

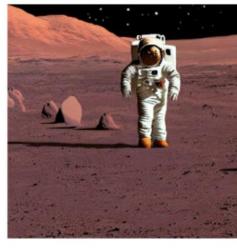
#### DALLE 2

#### Midjourney

#### n

#### Makin

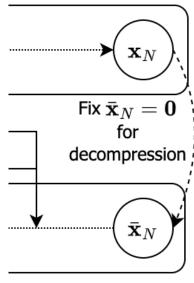
• During t  $x_N$  and iof their compres







Alone astronaut on Mars, mysterious, colorful, hyper realistic



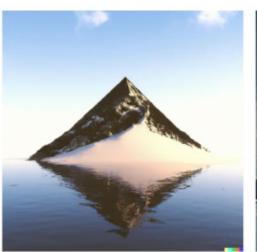
enoising

- Use you a row ar matchin from pu
- Yes, and these da

**Stable Diffusion** 



DALLE 2



Midjourney



large bus sitting next to a very tall



nk bed with a narrow shelf sitting underneath it.

Pyramid shaped mountain above a still lake, covered with snow

# Diffusion: a good entry point to continue your study of generative models

- **Stable diffusion:** Stable diffusion is a term used to describe a certain property of stochastic differential equations (SDEs) or diffusion processes.
- In the context of diffusion models, it refers to the stability of the diffusion process, ensuring that the SDE does not explode or diverge over time.
- Stable diffusion processes are crucial in designing diffusion models because they ensure that the generated samples remain within a reasonable range and do not become arbitrarily large or small.

# Diffusion: a good entry point to continue your study of generative models

- Diffusion models is a promising direction, but many more steps needed to get to a generative model like DALL-E or MidJourney!
- DALL-E and MidJourney models are not so to speak based on diffusion models, but have leveraged intuition from said models.
- They would be better described, for now, as generative models based on transformer architectures.
- Something for later in **Term 7 Computer Vision** course?

#### Conclusion

- 1. What is a **Stochastic Differential Equation**?
- 2. What are **typical uses** of SDEs?
- 3. Can I noise images progressively using SDEs?
- 4. Can you train a denoising autoencoder model that removes SDE noise from images?
- 5. What is the **diffusion process** and how to use it to train **diffusion models**?
- 6. What are the architectures of DALL-E and Midjourney models?

## Learn more about these topics

Out of class, for those of you who are curious

The one paper, usually cited for Diffusion models
 [Ho2020] J. Ho, X. Chen, A. Srinivas, Y. Duan, and P. Abbeel,
 "Denoising diffusion probabilistic models. Advances in Neural
 Information Processing Systems", 2020.

https://arxiv.org/abs/2006.11239

• [Dhariwal2021] P. Dhariwal, and A. Nichol, "Diffusion models beat GANs on ImageNet.", 2021.

https://arxiv.org/abs/2105.05233