

Probability Basics:
 $\mathbb{P}((X, Y) \in A) = \int_A p(x, y) \, dx \, dy.$
 $p(x) = \int_{-\infty}^{\infty} p(x, y) \, dy. \quad p(y|x) = \frac{p(x, y)}{p(x)}.$
 $p(x, y) = p(x)p(y|y|x).$ If X and Y are independent, then $p(x, y) = p(x)p(y).$
Likelihood of i.i.d. dataset: $p(\mathcal{D}) = \prod_{i=1}^n p(x_i)$ or $\log p(\mathcal{D}) = \sum_{i=1}^n \log p(x_i).$ where dataset $\mathcal{D} = \{x_1, \dots, x_n\}$ is i.i.d.
Bayes' theorem: $p(x|y) = \frac{p(x, y)}{p(y)} = \frac{p(y|x)p(x)}{p(y)}.$
Expectation:
 $\mathbb{E}[X] = \int_{-\infty}^{\infty} xp(x) \, dx, \quad \mathbb{E}[g(X)] = \int_{-\infty}^{\infty} g(x)p(x) \, dx,$
 $\mathbb{E}[X + Y] = \mathbb{E}[X] + \mathbb{E}[Y],$
 $\mathbb{E}[X | Y = y] = \int_{-\infty}^{\infty} xp(x | y) \, dx,$
 $\mathbb{E}[f(X)g(Y) | Y = y] = \mathbb{E}[f(X) | Y = y]g(y).$
Variance and Covariance:
 $\text{var}(X) = \mathbb{E}[(X - \mathbb{E}[X])^2] = \mathbb{E}[X^2] - \mathbb{E}[X]^2, \quad \text{cov}(X, Y) = \mathbb{E}[(X - \mathbb{E}[X])(Y - \mathbb{E}[Y])] = \mathbb{E}[XY] - \mathbb{E}[X]\mathbb{E}[Y],$
 $\text{var}(X + Y) = \text{var}(X) + \text{var}(Y) + 2\text{cov}(X, Y).$
Vector Variables: $\mathbf{x} = [x_1 \quad x_2 \quad \dots \quad x_n]^\top, \quad \mathbb{E}[\mathbf{x}] = [\mathbb{E}[x_1] \quad \mathbb{E}[x_2] \quad \dots \quad \mathbb{E}[x_n]]^\top$
If \mathbf{A} is a deterministic matrix, $\mathbb{E}[\mathbf{A}\mathbf{x}] = \mathbf{A}\mathbb{E}[\mathbf{x}].$
Covariance: $\text{cov}(\mathbf{x}) = \mathbf{\Sigma}_{xx} = \mathbb{E}[(\mathbf{x} - \mathbb{E}\mathbf{x})(\mathbf{x} - \mathbb{E}\mathbf{x})^\top],$
 $\text{cov}(\mathbf{A}\mathbf{x}) = \mathbf{A} \text{cov}(\mathbf{x}) \mathbf{A}^\top$
Cross-covariance: $\text{cov}(\mathbf{x}, \mathbf{y}) = \mathbf{\Sigma}_{xy} = \mathbb{E}[(\mathbf{x} - \mathbb{E}\mathbf{x})(\mathbf{y} - \mathbb{E}\mathbf{y})^\top]$
Matrix Calculus: $\frac{\partial(\mathbf{a}^\top \mathbf{x})}{\partial \mathbf{x}} = \mathbf{a}, \quad \frac{\partial(\mathbf{x}^\top \mathbf{A} \mathbf{x})}{\partial \mathbf{x}} = (\mathbf{A} + \mathbf{A}^\top) \mathbf{x}$
 $\frac{\partial(\mathbf{a}^\top \mathbf{x} \mathbf{b})}{\partial \mathbf{x}} = \mathbf{a} \mathbf{b}^\top, \quad \frac{\partial \det(\mathbf{X})}{\partial \mathbf{X}} = \det(\mathbf{X}) (\mathbf{X}^{-1})^\top,$
 $\frac{\partial(\mathbf{a}^\top \mathbf{X}^{-1} \mathbf{b})}{\partial \mathbf{X}} = -(\mathbf{X}^{-1})^\top \mathbf{a} \mathbf{b}^\top (\mathbf{X}^{-1})^\top$
Bayesian Inference:
Given a parametric model, the posterior is derived as:
 $p(\theta | \mathbf{x}) = \frac{p(\mathbf{x}|\theta)p(\theta)}{p(\mathbf{x})} \propto p(\mathbf{x} | \theta)p(\theta),$ where $p(\mathbf{x})$ is a normalization constant. If $p(\mathbf{x}) \propto f(\mathbf{x})$ for some function $f(\mathbf{x})$, then $p(\mathbf{x}) = cf(\mathbf{x})$ with $c = (\int f(\mathbf{x})d\mathbf{x})^{-1}.$
Conjugate Distributions:
If prior and posterior share the same form, they are conjugate:
 $p(\theta | \mathbf{x}) \propto p(\mathbf{x} | \theta)p(\theta).$ The prior $p(\theta)$ is called the conjugate prior of the likelihood $p(\mathbf{x} | \theta).$
1. Allows for analytical closed form solutions and easy to interpret.
2. May lack flexibility to complex data, requiring MCMC.
Conjugate Prior for Binomial:
Given $s \sim \text{Bin}(n, \theta)$ and prior $\theta \sim \text{Beta}(a, b),$ the posterior is:
 $p(\theta | s) \propto \text{Bin}(s | \theta, n) \cdot \text{Beta}(\theta | a, b)$
 $\propto \theta^s (1 - \theta)^{n-s} \cdot \theta^{a-1} (1 - \theta)^{b-1}$
 $\propto \theta^{s+a-1} (1 - \theta)^{n-s+b-1}.$ Thus, the posterior follows:
 $p(\theta | s) = \text{Beta}(\theta | s + a, n - s + b)$
Categorical Distribution:
• A categorical variable X follows: $\text{Cat}(x | \theta_1, \dots, \theta_K)$ with parameters $\theta_k \geq 0, \sum_{k=1}^K \theta_k = 1; \text{Cat}(x | \theta_1, \dots, \theta_K) = \theta_x.$
• Given i.i.d. samples $X_i \sim \text{Cat}(\theta_1, \dots, \theta_K),$ the joint probability of $\mathcal{D} = \{X_1, \dots, X_n\}$ is: $p(\mathcal{D}) = \prod_{i=1}^n \text{Cat}(x_i | \theta_1, \dots, \theta_K) = \prod_{i=1}^n \prod_{k=1}^K \theta_k^{1\{x_i=k\}}.$ Using count notation $N_k = \sum_{i=1}^n 1\{x_i = k\},$ we get: $p(\mathcal{D}) = \prod_{k=1}^K \theta_k^{N_k}.$
Gaussian (Normal) Distribution:
A random variable X follows a normal distribution:
 $\mathcal{N}(x | \mu, \sigma^2) = \frac{1}{\sqrt{2\pi\sigma^2}} e^{-\frac{1}{2\sigma^2}(x-\mu)^2}.$
If $X \sim \mathcal{N}(\mu, \sigma^2),$ then: $aX \sim \mathcal{N}(a\mu, a^2\sigma^2)$ for any $a \in \mathbb{R}$
 $X + c \sim \mathcal{N}(\mu + c, \sigma^2).$ If $Z \sim \mathcal{N}(0, 1),$ then $X = \sigma Z + \mu \sim \mathcal{N}(\mu, \sigma^2)$ and $Y \sim \mathcal{N}(\xi, \nu^2)$ are independent: $X + Y \sim \mathcal{N}(\mu + \xi, \sigma^2 + \nu^2).$
Multivariate Gaussian Distribution:
A random vector \mathbf{X} follows a multivariate normal distribution:
 $\mathcal{N}(\mathbf{x} | \mu, \Sigma) = \frac{1}{(2\pi)^K/2(\det \Sigma)^{1/2}} \exp\left(-\frac{1}{2}(\mathbf{x} - \mu)^\top \Sigma^{-1}(\mathbf{x} - \mu)\right).$
• $\mathbf{X} = [X_1, \dots, X_n]^\top$ is jointly Gaussian if for any vector $\mathbf{a} \in \mathbb{R}^n,$ the linear combination: $\mathbf{a}^\top \mathbf{X} = \sum_{i=1}^n a_i X_i$ is Gaussian.
• If $\mathbf{Z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I}),$ then: $\mathbf{X} = \mathbf{A}\mathbf{Z} + \mu$ is jointly Gaussian with mean and covariance: $\mathbb{E}[\mathbf{X}] = \mu, \quad \text{cov}(\mathbf{X}) = \mathbf{A}\mathbf{A}^\top.$
• $\mu = \begin{bmatrix} \mu_X \\ \mu_Y \end{bmatrix}, \quad \Sigma = \begin{bmatrix} \sigma_X^2 & \text{cov}(X, Y) \\ \text{cov}(X, Y) & \sigma_Y^2 \end{bmatrix}.$
 $\rho_{X,Y} = \frac{\text{cov}(X,Y)}{\sigma_X \sigma_Y}.$
Likelihood Functions for Common Distributions:
1. Bernoulli: If $x_i \sim \text{Bern}(\theta)$ i.i.d., where $\theta \in [0, 1],$ then:
 $p(\mathcal{D} | \theta) = \theta^{N_1} (1 - \theta)^{N_0}, \quad N_k = \sum_{i=1}^n 1\{x_i = k\}.$
2. Exponential: If $x_i \sim \text{Exp}(\lambda)$ i.i.d., where $\lambda > 0,$ then:
 $p(\mathcal{D} | \lambda) = \lambda^n \exp(-\lambda \sum_{i=1}^n x_i).$
3. Gaussian: If $x_i \sim \mathcal{N}(\mu, \sigma^2)$ i.i.d., then $\theta = (\mu, \sigma^2)$ and:
 $p(\mathcal{D} | \theta) = \frac{1}{(2\pi)^{\frac{n}{2}} \sigma^n} \exp\left(-\frac{1}{2\sigma^2} \sum_{i=1}^n (x_i - \mu)^2\right).$
Maximum Likelihood Estimation (MLE):
The MLE of parameter θ is: $\theta_{\text{ML}} = \arg \max_{\theta} p(\mathcal{D} | \theta).$ For i.i.d. data $\mathcal{D} = \{x_1, \dots, x_n\},$ we maximize the log-likelihood:
 $\log p(\mathcal{D} | \theta) = \sum_{i=1}^n \log p(x_i | \theta).$
Bernoulli MLE: If $x_i \sim \text{Bern}(\theta),$ then:
 $p(\mathcal{D} | \theta) = \theta^{N_1} (1 - \theta)^{N_0},$
 $\log p(\mathcal{D} | \theta) = N_1 \log \theta + N_0 \log(1 - \theta).$ Solving $\frac{\partial}{\partial \theta} \log p(\mathcal{D} | \theta) = 0:$ $\theta_{\text{ML}} = \frac{N_1}{N_0 + N_1} = \frac{N_1}{n}.$
Exponential MLE: If $x_i \sim \text{Exp}(\lambda),$ then:
 $\log p(\mathcal{D} | \lambda) = n \log \lambda - \lambda \sum_{i=1}^n x_i.$ Solving $\frac{\partial}{\partial \lambda} \log p(\mathcal{D} | \lambda) = 0:$ $\lambda_{\text{ML}} = \frac{n}{\sum_{i=1}^n x_i}.$
Linear Regression Model:
 $\mathbf{x} = (x_1, x_2, \dots, x_D), y = \mathbf{w}^\top \mathbf{x} + \epsilon, \epsilon \sim \mathcal{N}(0, \sigma^2).$
Thus, the likelihood is: $p(y | \mathbf{x}, \mathbf{w}) = \mathcal{N}(y | \mathbf{w}^\top \mathbf{x}, \sigma^2).$
Basis Function Expansion: model non-linear relationships using

basis functions: $\phi(\mathbf{x}) = [\phi_1(\mathbf{x}) \quad \phi_2(\mathbf{x}) \quad \dots \quad \phi_M(\mathbf{x})]^\top$
 $y = \mathbf{w}^\top \phi(\mathbf{x}) + \epsilon;$ For $\mathbf{x} = \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}, \phi(\mathbf{x}) = [1, x_1, x_2, x_1^2, x_2^2]^\top,$
 $y = w_1 + w_2 x_1 + w_3 x_2 + w_4 x_1^2 + w_5 x_2^2 + \epsilon.$
MLE for \mathbf{w} : Given i.i.d. training data $\mathcal{D} = \{(\mathbf{x}_i, y_i)\},$
 $\log p(\mathbf{y} | \Phi, \mathbf{w}) = -\frac{1}{2\sigma^2} \|\Phi \mathbf{w} - \mathbf{y}\|^2 + \text{const},$ where $\Phi = [\mathbf{x}_1 \dots \mathbf{x}_n]^\top$ is the design matrix. Maximizing w.r.t. \mathbf{w} gives the least squares solution: $\mathbf{w}_{\text{ML}} = (\Phi^\top \Phi)^{-1} \Phi^\top \mathbf{y}.$
Model Evaluation Metrics:
• $\text{RSS} = \sum_{i=1}^n (y_i - \hat{y}_i)^2. \quad \bullet \text{RMSE} = \sqrt{\frac{1}{n} \text{RSS}}.$
• $R^2 = 1 - \frac{\text{RSS}}{\text{TSS}} = 1 - \frac{\text{RSS}}{\sum_{i=1}^n (y_i - \bar{y})^2}.$
Maximum A Posteriori (MAP) Estimation:
• The MAP estimate maximizes the posterior:
 $\theta_{\text{MAP}} = \arg \max_{\theta} p(\theta | \mathcal{D}).$ Using Bayes' theorem:
 $p(\theta | \mathcal{D}) \propto p(\mathcal{D} | \theta)p(\theta),$
 $\log p(\theta | \mathcal{D}) = \log p(\mathcal{D} | \theta) + \log p(\theta) + \text{const}.$
• Given: $p(\theta | \mathcal{D}) = \theta^{N_1} (1 - \theta)^{N_0}, \quad p(\theta) = \text{Beta}(\theta | a, b),$
 $\log p(\mathcal{D} | \theta)p(\theta) = (N_1 + a - 1) \log \theta + (N_0 + b - 1) \log(1 - \theta).$ Solving $\frac{\partial}{\partial \theta} \log p(\mathcal{D} | \theta) = 0:$ $\theta_{\text{MAP}} = \frac{N_1 + a - 1}{n + a + b - 2}.$
Classification and Naïve Bayes:
• Classification Rule: Given feature vector \mathbf{x} and class label $y \in \{1, \dots, K\}, \delta(\mathbf{x}) = k$ if $p(y = k | \mathbf{x})$ is maximized.
• Naïve Bayes Classifier: Assuming conditional independence, $p(\mathbf{x} | y = c, \theta) = \prod_{d=1}^D p(x_d | \theta_{dc}).$ Using Bayes' rule:
 $p(y = c | \mathbf{x}, \theta) \propto \pi(c) \prod_{d=1}^D p(x_d | \theta_{dc}),$ where $\pi(c)$ is the prior probability of class $c.$
Mixture Models and Gaussian Mixture Model (GMM):
• Suppose an observation \mathbf{x} can be generated from one of K possible probability density functions (pdfs): $p(\mathbf{x} | \boldsymbol{\eta}_1), \dots, p(\mathbf{x} | \boldsymbol{\eta}_K).$ The generating index z follows a categorical distribution:
 $p(z) = \text{Cat}(z | \boldsymbol{\pi}).$ Since z is unobserved, it is a **latent variable**. The marginal distribution is: $p(\mathbf{x} | \boldsymbol{\theta}) = \sum_{k=1}^K \pi(k) p(\mathbf{x} | \boldsymbol{\eta}_k),$ where $\boldsymbol{\theta} = (\boldsymbol{\pi}, \{\boldsymbol{\eta}_k\}_{k=1}^K).$
• GMM: $p(\mathbf{x} | \boldsymbol{\theta}) = \sum_{k=1}^K \pi(k) \mathcal{N}(\mathbf{x} | \boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k),$ where $\boldsymbol{\theta} = (\boldsymbol{\pi}, \{\boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k\}_{k=1}^K).$ Given i.i.d. observations $\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n,$ $p(\mathbf{x}_i | \boldsymbol{\theta}) = \sum_{k=1}^K \pi(k) p(\mathbf{x}_i | \boldsymbol{\eta}_k).$
 $\log p(\mathbf{x}_1, \dots, \mathbf{x}_n | \boldsymbol{\theta}) = \sum_{i=1}^n \log \sum_{k=1}^K \pi(k) p(\mathbf{x}_i | \boldsymbol{\eta}_k).$
Algorithmic Issues in Mixture Models:
• Singularity: If for some $k,$ we set $\boldsymbol{\mu}_k = \mathbf{x}_i$ and $\sigma_k \rightarrow 0,$ then: $\mathcal{N}(\mathbf{x}_i | \boldsymbol{\mu}_k, \sigma_k \mathbf{I}) \propto \frac{1}{\sigma_k} \rightarrow \infty.$
• Unidentifiability: no unique global optimum for log-likelihood function.
• Optimization Challenges: non-convex, hard to solve.
• If latent variables z_1, \dots, z_n are observed, the likelihood simplifies to: $\log p(\mathbf{x}_1, z_1, \dots, \mathbf{x}_n, z_n | \boldsymbol{\theta}) = \sum_{i=1}^n (\log \pi[z_i] + \log p(\mathbf{x}_i | \boldsymbol{\eta}_{z_i})).$ much easier to maximize.
Gaussian Mixture Model (GMM):
• Observed data $\mathbf{x}_1, \dots, \mathbf{x}_n \in \mathbb{R}^D$ are generated from a mixture of K Gaussian distributions:
 $p(\mathbf{x} | \boldsymbol{\theta}) = \sum_{k=1}^K \pi(k) \mathcal{N}(\mathbf{x} | \boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k)$
 $\boldsymbol{\theta} = (\pi(k), \boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k)_{k=1}^K.$
• Complete data likelihood: $\log p(\mathbf{y}_1, \dots, \mathbf{y}_n | \boldsymbol{\theta}) = \sum_{k=1}^K \sum_{i: z_i=k} (\log \pi(k) + \log \mathcal{N}(\mathbf{x}_i | \boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k)).$
• The MLE estimates for $\boldsymbol{\mu}_k$ and $\boldsymbol{\Sigma}_k:$ $\hat{\boldsymbol{\mu}}_k = \frac{1}{n} \sum_{i: z_i=k} \mathbf{x}_i,$
 $\hat{\boldsymbol{\Sigma}}_k = \frac{1}{n} \sum_{i: z_i=k} (\mathbf{x}_i - \hat{\boldsymbol{\mu}}_k)(\mathbf{x}_i - \hat{\boldsymbol{\mu}}_k)^\top.$
Expectation-Maximization (EM) Algorithm:
Complete data \mathbf{y} can't be known directly. Still since $\log p(\mathbf{x} | \boldsymbol{\theta})$ is hard to optimize. Try maximize the expectation
 $\mathbb{E}_{p(\mathbf{y}|\mathbf{x}, \boldsymbol{\theta})} [\log p(\mathbf{y} | \boldsymbol{\theta}) | \mathbf{x}, \hat{\boldsymbol{\theta}}].$
EM Steps: 1. Initialize $\boldsymbol{\theta}^{(0)}.$ 2. **E-step:** Compute $Q(\boldsymbol{\theta} | \boldsymbol{\theta}^{(m)}) = \mathbb{E}_{p(\mathbf{y}|\mathbf{x}, \boldsymbol{\theta}^{(m)})} [\log p(\mathbf{y} | \boldsymbol{\theta}) | \mathbf{x}, \boldsymbol{\theta}^{(m)}] = \int \log p(\mathbf{y} | \boldsymbol{\theta}) p(\mathbf{y} | \mathbf{x}, \boldsymbol{\theta}^{(m)}) d\mathbf{y}.$
3. M-step: Update $\boldsymbol{\theta}$ by $\boldsymbol{\theta}^{(m+1)} = \arg \max_{\boldsymbol{\theta} \in \Theta} Q(\boldsymbol{\theta} | \boldsymbol{\theta}^{(m)}).$
4. Repeat until convergence.
EM for GMM:
Given $\pi^{(0)}(k), \boldsymbol{\mu}_k^{(0)}, \boldsymbol{\Sigma}_k^{(0)}$ for $k = 1, \dots, K.$
 $L^{(0)} = \frac{1}{n} \sum_{i=1}^n \log \left(\sum_{k=1}^K \pi^{(0)}(k) \mathcal{N}(\mathbf{x}_i | \boldsymbol{\mu}_k^{(0)}, \boldsymbol{\Sigma}_k^{(0)}) \right).$
Repeat:
1. **E-step:** Compute responsibilities:
 $r_{ik}^{(m)} = \frac{\pi^{(m)}(k) \mathcal{N}(\mathbf{x}_i | \boldsymbol{\mu}_k^{(m)}, \boldsymbol{\Sigma}_k^{(m)})}{\sum_{k'} \pi^{(m)}(k') \mathcal{N}(\mathbf{x}_i | \boldsymbol{\mu}_{k'}^{(m)}, \boldsymbol{\Sigma}_{k'}^{(m)})}, n_k^{(m)} = \sum_{i=1}^n r_{ik}^{(m)}.$
2. **M-step:** Update parameters: $\pi^{(m+1)}(k) = \frac{n_k^{(m)}}{n},$
 $\boldsymbol{\mu}_k^{(m+1)} = \frac{1}{n_k^{(m)}} \sum_{i=1}^n r_{ik}^{(m)} \mathbf{x}_i, \quad \boldsymbol{\Sigma}_k^{(m+1)} = \frac{1}{n_k^{(m)}} \sum_{i=1}^n r_{ik}^{(m)} (\mathbf{x}_i - \boldsymbol{\mu}_k^{(m+1)})(\mathbf{x}_i - \boldsymbol{\mu}_k^{(m+1)})^\top.$
3. **Compute:** $L^{(m+1)} = \frac{1}{n} \sum_{i=1}^n \log \left(\sum_{k=1}^K \pi^{(m+1)}(k) \mathcal{N}(\mathbf{x}_i | \boldsymbol{\mu}_k^{(m+1)}, \boldsymbol{\Sigma}_k^{(m+1)}) \right).$
4. **untill:** $|L^{(m+1)} - L^{(m)}| \leq \epsilon.$
K-Means Algorithm
GMM with $\Sigma_k = \sigma^2 I$ and $\pi(k) = 1/K$ are fixed. Only μ_k are inferred. 1. **E-step:** Assign each x_i to its nearest cluster center: $k_i = \arg \min_k \|x_i - \mu_k^{(m)}\|^2.$ Define hard assignment:
 $r_{ik}^{(m)} = 1$ if $k = k_i$, otherwise $r_{ik}^{(m)} = 0.$
 $Q(\boldsymbol{\theta} | \boldsymbol{\theta}^{(m)}) = -\frac{1}{2\sigma^2} \sum_{i=1}^n \|\mathbf{x}_i - \mu_{k_i}^{(m)}\|^2 + \text{const}.$
2. **M-step:** Update cluster centers: $\mu_k^{(m+1)} = \frac{1}{N_k} \sum_{i: k_i=k} x_i,$ where $N_k = \sum_{i=1}^n 1\{k_i = k\}.$
3. Update each cluster center using the mean of its assigned points.
EM Algorithm for MAP Estimation

Given data \mathbf{x} , the posterior is $p(\theta | \mathbf{x}) = \frac{p(\mathbf{x}|\theta)p(\theta)}{p(\mathbf{x})}$
estimate is $\theta_{MAP} = \arg \max_{\theta} (\log p(\mathbf{x} | \theta) + \log p(\theta)).$
Algorithm:
1. Pick initial guess $\theta^{(0)}.$
2. **E-step:** At iteration $m + 1,$ compute $Q(\theta | \theta^{(m)}) = \int p(\mathbf{y} | \mathbf{x}, \theta^{(m)}) \log p(\mathbf{y} | \theta) d\mathbf{y}.$
3. **M-step:** Update $\theta^{(m+1)} = \arg \max_{\theta} \left(Q(\theta | \theta^{(m)}) + \log p(\theta) \right).$
4. Repeat until convergence.
Transition Matrix
• Each row of \mathbf{T} sums to one:
 $\sum_{j=1}^M T(i, j) = \sum_{j=1}^M p_{x[t]|x[t-1]}(j | i) = 1.$ \mathbf{T} is a (row) stochastic matrix.
• Given initial distribution $\mathbf{p}_0 = [p_0(1), p_0(2), \dots, p_0(M)],$
 $\mathbf{p}_1 = \mathbf{p}_0 \mathbf{T}.$ For general $t,$ let $\mathbf{p}_t = [p_t(1), \dots, p_t(M)],$ then: $\mathbf{p}_t = \mathbf{p}_0 \mathbf{T}^t.$
MLE for Transition Matrix
• Estimate prior π and transition matrix \mathbf{T} from training data:
 $p(\mathbf{x}[0], \dots, \mathbf{x}[t] | \pi, \mathbf{T}) = \pi(\mathbf{x}[0]) \mathbf{T}(\mathbf{x}[0], \mathbf{x}[1]) \dots \mathbf{T}(\mathbf{x}[t-1], \mathbf{x}[t]).$ Given n observed sequences $\mathcal{D} = \{\mathbf{x}_1[0 : t_1], \dots, \mathbf{x}_n[0 : t_n]\},$ each of varying length $t_i + 1,$ assume all data points follow the same $\mathbf{T}.$
• $\log p(\mathcal{D} | \pi, \mathbf{T}) = \sum_{i=1}^n \log \pi(\mathbf{x}_i[0]) + \sum_{i=1}^n \sum_{t=1}^{t_i} \log \mathbf{T}(\mathbf{x}_i[t-1], \mathbf{x}_i[t])$
 $= \sum_{x=1}^M N_x \log \pi(x) + \sum_{x=1}^M \sum_{y=1}^M N_{xy} \log \mathbf{T}(x, y).$
 $N_x = \sum_{i=1}^n 1(\mathbf{x}_i[0] = x),$
 $N_{xy} = \sum_{i=1}^n \sum_{t=1}^{t_i} 1(\mathbf{x}_i[t-1] = x, \mathbf{x}_i[t] = y).$
• $\hat{\pi}(x) = \frac{N_x}{n}, \hat{\mathbf{T}}(x, y) = \frac{N_{xy}}{\sum_{z=1}^M N_{xz}}.$
• May predict certain strings are impossible if data have zero count of certain states (this is a form of overfitting).
HMM: Hidden Markov Model
A hidden Markov model consists of:
• A discrete state Markov chain with **hidden states** or latent variables $z[t] \in \{1, \dots, M\}, t = 0, 1, \dots,$ with initial pdf π and transition matrix $\mathbf{T}.$
• An observation model with emission probabilities $p(\mathbf{x}[t] | z[t]) = p(\mathbf{x}[t] | \phi_{z[t]}),$ where $\phi = (\phi_1, \dots, \phi_M).$ Applications: (1) Long-range dependencies; (2) Speech recognition; (3) Gene finding; (4) Emission probabilities (Gaussian example).
Baum-Welch Algorithm for HMM Training
1. Initialize $\theta^{(0)}.$
2. **E step:** Use Forward-Backward Algorithm to compute $\gamma_{i,t}(z) = p(z_i[t] = z | \mathbf{x}_i[0 : t_i], \theta^{(m)}) \propto \alpha_z(z) \beta_z(z).$
 $\xi_{i,t}(z, z') = p(z_i[t-1] = z, z_i[t] = z' | \mathbf{x}_i[0 : t_i], \theta^{(m)}) \propto \alpha_{t-1}(z) p(\mathbf{x}_i[t] | z_i[t]) \beta_{t'}(z') p(z_i[t] = z' | z_i[t-1] = z).$
3. **M step:** Update parameters $\hat{\pi}(z) = \frac{\sum_{i=1}^n \gamma_{i,0}(z)}{n},$
 $\hat{T}(z, z') = \frac{\sum_{i=1}^n \sum_{t=1}^{t_i} \xi_{i,t}(z, z')}{\sum_u \sum_{i=1}^n \sum_{t=1}^{t_i} \xi_{i,t}(z, u)}, \hat{\phi}_z = \text{emission probability parameters}.$
4. Repeat E and M steps.
Inference in HMMs
• Filtering: Estimate latent state $p(z[t] | \mathbf{x}[0 : t])$ using observations up to time t (Forward Algorithm).
• Smoothing: Estimate $p(z[t] | \mathbf{x}[0 : T]),$ using both past and future observations (Forward-Backward Algorithm).
• Fixed-lag smoothing: Estimate $p(z[t-l] | \mathbf{x}[0 : t])$ for online inference (Forward-Backward Algorithm).
• Prediction: Estimate $p(z[t+h] | \mathbf{x}[0 : t]),$ where $h > 0$ is the prediction horizon: $p(z[t+h] | \mathbf{x}[0 : t]) = \sum_{z[t], \dots, z[t+h-1]} p(z[t+h] | z[t+h-1]) p(z[t+h-1] | z[t+h-2]) \dots p(z[t+1] | z[t]) \cdot p(z[t] | \mathbf{x}[0 : t]).$
• MAP sequence: using Viterbi Algorithm, most probable sequence $\mathbf{z}^*[0 : T] = \arg \max_{\mathbf{z}[0:T]} p(z[0 : T] | \mathbf{x}[0 : T]).$
Sampling Using Cdf
If $X \sim F,$ then: $\mathbb{P}(X \leq x) = F(x), \quad \text{Let } U \sim \text{Unif}(0, 1),$ define: $X = F^{-1}(U) \Rightarrow X \sim F.$
 $\mathbb{P}(F^{-1}(U) \leq x) = \mathbb{P}(U \leq F(x)) = F(x).$
Transformations
If $Y = f(X),$ then $p_Y(y) = \sum_{k=1}^K \frac{p_X(x_k)}{|f'(x_k)|},$ where x_k are the solutions to $f(x) = y$ **Note:** Requires solving $f(x) = y$ and knowing $f'(x)$
Rejection Sampling
• $p(z) = \frac{1}{M} \tilde{p}(z),$ where M unknown. Choose proposal $q(z),$ and constant $k \geq \frac{\tilde{p}(z)}{q(z)}$ for all $z.$
• $\text{supp}(p) \subseteq \text{supp}(q),$ where $\text{supp } p = \{z : p(z) > 0\}.$
• **Sampling Procedure:** 1. Sample $z \sim q(z)$ 2. Sample $u \sim \text{Unif}(0, kq(z))$ 3. Accept z if $u \leq \tilde{p}(z).$
• $\mathbb{P}(z \text{ accepted}) = \int \frac{\tilde{p}(z)}{kq(z)} q(z) dz = \frac{M}{k}.$
• Choose smallest possible k s.t. $kq(z) \geq \tilde{p}(z) \forall z.$
• Accepted $z \sim p(z)$
 $\mathbb{P}(z \leq z_0 | \text{accepted}) = \frac{1}{M} \int_{z \leq z_0} \tilde{p}(z) dz \rightarrow$ Accepted samples follow the CDF of $p(z)$
Importance Sampling
• Estimate expectation $\mathbb{E}_p[f(z)] = \int f(z)p(z) dz,$ where $p(z)$ is hard to sample.
• Use proposal $q(z)$ and importance weights $w(z) = \frac{p(z)}{q(z)},$ rewrite as $\mathbb{E}_p[f(z)] = \int f(z)w(z)q(z) dz \approx \frac{1}{n} \sum_{i=1}^n w(z_i) f(z_i),$ where $z_i \sim q(z).$
• $\text{supp}(f(\cdot)p(\cdot)) \subseteq \text{supp}(q).$ Keep all samples, no need $q(z) \geq p(z).$ Better matches well, if $q(z)$ is large where $|f(z)|p(z)$ is large.
• If only know how to compute $p(z)$ and $q(z)$ up to normalizing constants, use normalized weights:
 $\mathbb{E}_p[f(z)] \approx \sum_{i=1}^n w_n(z_i) f(z_i), w_n(z_i) = \frac{w(z_i)}{\sum_j w(z_j)}.$
Sampling Importance Resampling (SIR)
• Convert importance weighted samples into unweighted samples

• Steps:
1. Sample z_1, \dots, z_n from $q(z)$.
2. Compute weights $w_n(z_1), \dots, w_n(z_n)$.
3. Resample with replacement from $\{z_1, \dots, z_n\}$ using weights $(w_n(z_1), \dots, w_n(z_n))$.
• Each sample \tilde{z}_i drawn from multinomial over $\{z_1, \dots, z_n\}$ with weights $w_n(z_i)$. asymptotically for large $n \rightarrow \infty$,
 $\mathbb{P}(\tilde{z} \leq a) = \sum_{i=1}^n w_n(z_i) \mathbb{1}_{\{z_i \leq a\}} \rightarrow \int_{z \leq a} p(z) dz$.
SIR for Bayesian Inference
• Take unnormalized $\tilde{p}(\theta) = p(\mathcal{D} \mid \theta)p(\theta)$, sample $\theta_1, \dots, \theta_n$ from $q(\theta) = p(\theta)$.
• $w_n(\theta_i) = \frac{\tilde{p}(\theta_i)/q(\theta_i)}{\sum_j \tilde{p}(\theta_j)/q(\theta_j)} = \frac{p(\mathcal{D}|\theta_i)}{\sum_j p(\mathcal{D}|\theta_j)}$.
• Resampling $\theta_1, \dots, \theta_n$ according to weights $(w_n(\theta_1), \dots, w_n(\theta_n))$.

Sampling for EM
• Observed incomplete data \mathbf{x} , while complete is (\mathbf{x}, \mathbf{z}) . Need to compute: $Q(\theta \mid \theta^{(m)}) = \int p(\mathbf{z} \mid \mathbf{x}, \theta^{(m)}) \log p(\mathbf{x}, \mathbf{z} \mid \theta) dz$.
• $Q(\theta \mid \theta^{(m)}) \approx \frac{1}{n} \sum_{i=1}^n \log p(\mathbf{x}, \mathbf{z}_i \mid \theta)$, where $\mathbf{z}_i \sim p(\mathbf{z} \mid \mathbf{x}, \theta^{(m)})$.
• Rejection and importance sampling not suitable for high-dimensional $\mathbf{z} \rightarrow$ need MCMC methods

Stationary Distribution
Consider homogeneous Markov chain with transition probability $p(x_t = y \mid x_{t-1} = x) = \mathbf{T}(x, y)$. π is a stationary distribution if $\sum_x \pi(x) \mathbf{T}(x, y) = \pi(y)$ for all states y . Also called invariant distribution — does not change over time in the chain. If M states, \mathbf{T} is an $M \times M$ matrix with $\pi \mathbf{T} = \pi$.
Asymptotic Steady State
Initial distribution: π_0 . Markov evolution:
 $\pi_1 = \pi_0 \mathbf{T}$, $\pi_2 = \pi_1 \mathbf{T}$, \dots , $\pi_k = \pi_0 \mathbf{T}^k$. If the limit $\pi = \lim_{k \rightarrow \infty} \pi_k = \lim_{k \rightarrow \infty} \pi_0 \mathbf{T}^k$ exist, π must satisfy $\pi \mathbf{T} = \pi \rightarrow$ a stationary distribution

Reversible Markov Chain (MC)
• Sufficient condition but not necessary condition for π to be stationary: $\pi(x) \mathbf{T}(x, y) = \pi(y) \mathbf{T}(y, x)$ for all $x, y \rightarrow$ This chain is *reversible*. Summing both sides over x :
 $\sum_x \pi(x) \mathbf{T}(x, y) = \sum_x \pi(x) \mathbf{T}(y, x) = \pi(y) \sum_x \mathbf{T}(y, x) = \pi(y)$.
• Sampling from a distribution $\pi(x)$: Design transition $\mathbf{T}(x, y)$ such that (1) All the ergodicity conditions hold (2) Usually make \mathbf{T} reversible and aperiodic (3) $\pi(x)$ is stationary distribution.
• MCMC basis: Generate sample path for the Markov chain z_0, z_1, \dots, z_n starting from any z_0 . If n large, then $p(z_n) \approx \pi(z_n) \rightarrow$ obtain samples from π
Metropolis-Hastings Algorithm

• Sample from $\pi(\mathbf{x})$ (e.g., $\mathcal{X} = \mathbb{R}^{1000}$). Assume we can compute unnormalized density $\tilde{\pi}(\mathbf{x})$. Choose transition probability $q(\mathbf{x}, \mathbf{y})$ that is irreducible, aperiodic, and easy to sample — this is the *proposal distribution*.
• Let Z_0, Z_1, \dots be the chain states. At step m :
1. Let $\mathbf{x} = Z_{m-1}$. 2. Sample $\mathbf{y} \sim q(\mathbf{x}, \cdot)$.
3. Accept \mathbf{y} with probability $A(\mathbf{x}, \mathbf{y}) = \min\left(1, \frac{\tilde{\pi}(\mathbf{y})q(\mathbf{y}, \mathbf{x})}{\tilde{\pi}(\mathbf{x})q(\mathbf{x}, \mathbf{y})}\right)$.
4. If accepted, $Z_m = \mathbf{y}$; else $Z_m = Z_{m-1}$.
• (1) Sequence Z_0, Z_1, \dots is a Markov chain (only depends on Z_{m-1}) (2) Like rejection/importance sampling, does not need normalization constant of $\tilde{\pi}$ (3) Suitable for high-dimensional $\pi(\mathbf{x})$ since sampling is from simple $q(\mathbf{x}, \cdot)$

Proposal Distributions
• MH chain Z_0, Z_1, \dots is designed to converge to stationary distribution $\pi(\cdot)$. For large m , $Z_m \sim \pi$ approximately.
• Burn-in period: discard first 1000 to 5000 samples
• Choice of Proposal Distribution:
1. $q(\mathbf{x}, \mathbf{y}) = q(\mathbf{y} - \mathbf{x})$ — *random walk MH*.
(1) $\mathbf{y} - \mathbf{x} \sim \mathcal{N}(0, \Sigma)$ — Gaussian centered at \mathbf{x} .
(2) $\mathbf{y} - \mathbf{x} \sim \text{Unif}[-\delta, \delta]^d$ — Uniform around \mathbf{x} .
• (3) If $q(\mathbf{x}, \mathbf{y}) = q(\mathbf{y}, \mathbf{x})$, then $A(\mathbf{x}, \mathbf{y}) = \min\left(1, \frac{\tilde{\pi}(\mathbf{y})}{\tilde{\pi}(\mathbf{x})}\right)$ known as the Metropolis Algorithm.
(4) Variance affects — small \Rightarrow slow, high \Rightarrow high rejection rates.
2. Independence Chain MH: $q(\mathbf{x}, \mathbf{y}) = q(\mathbf{y})$, i.e. next state is independent of current. Works well if $q(\mathbf{y})$ closely approximates $\pi(\mathbf{y})$ and is heavy-tailed
3. Exploiting Structure of π : Suppose $\pi(\mathbf{x}) \propto \psi(\mathbf{x})h(\mathbf{x})$, with $h(\mathbf{x})$ is a density that can be sampled from, bounded function $\psi(\mathbf{x})$. Choose $q(\mathbf{x}, \mathbf{y}) = h(\mathbf{y})$, then:

$A(\mathbf{x}, \mathbf{y}) = \min\left(1, \frac{\tilde{\pi}(\mathbf{y})q(\mathbf{y}, \mathbf{x})}{\tilde{\pi}(\mathbf{x})q(\mathbf{x}, \mathbf{y})}\right) = \min\left(1, \frac{\psi(\mathbf{y})h(\mathbf{y})h(\mathbf{x})}{\psi(\mathbf{x})h(\mathbf{x})h(\mathbf{y})}\right) = \min\left(1, \frac{\psi(\mathbf{y})}{\psi(\mathbf{x})}\right)$
Proposal Variance • Important to tune proposal variance σ .
• σ too small \rightarrow slow convergence, high acceptance rate, takes a long time to explore the whole space.
• σ too large \rightarrow big steps, low acceptance rate, stuck for long time.
• Rules of thumb: (1) Random walk MH: target acceptance rate of 0.25 to 0.5. (2) Independence chain MH: acceptance rate close to 1
Burn-In Discard early samples before chain reaches stationary distribution. Hard to know when. E.g. $x_0 \sim \text{Unif}(\{0, 1, \dots, 20\})$, takes over 400 steps to “forget” initial state.

Thinning Break dependencies between samples by taking every \bar{d} -th sample. Useful when σ is too large \rightarrow MC stuck for long time at same location. Subsample every \bar{d} samples: z_0, z_d, z_{2d}, \dots

Gibbs Sampling
Special case of Metropolis-Hastings, for multivariate $p(z_1, \dots, z_d)$ (hard to sample jointly when d is large). For each i , define $\mathbf{z}_{-i} = \{z_1, \dots, z_{i-1}, z_{i+1}, \dots, z_d\}$. If we can compute full conditionals $p(z_i \mid \mathbf{z}_{-i})$, we can perform Gibbs sampling:

• Initialize $(z_1^{(0)}, \dots, z_d^{(0)})$.
• For each iteration k , sequentially sample:
 $z_1^{(k)} \sim p(\cdot \mid z_2^{(k-1)}, \dots, z_d^{(k-1)})$
 $z_2^{(k)} \sim p(\cdot \mid z_1^{(k)}, z_3^{(k-1)}, \dots)$...
 $z_j^{(k)} \sim p(\cdot \mid z_1^{(k)}, \dots, z_{j-1}^{(k)}, z_{j+1}^{(k-1)}, \dots)$...

from $p(z)$.
 $\mathbf{z}^{(k)} \sim p(\cdot \mid \mathbf{z}^{(k)}_1, \dots, \mathbf{z}^{(k)}_{d-1})$.
• Discard burn-in samples.
• Generating Approximate i.i.d. Samples: (1) Run r independent Gibbs chains of length m , use final sample from each sequence. (2) Run one long sequence, discard burn-in, take every d -th sample.
• Getting conditionals distributions:
 $p(z_1 \mid z_2, \dots, z_d) = \frac{p(z_1, \dots, z_d)}{p(z_2, \dots, z_d)}$ (1) Start with joint $p(z_1, \dots, z_d)$. (2) Drop constants not depend on z_1 . (3) Use knowledge of well-known distributions to find the distribution $p(z_1 \mid z_2, \dots, z_d)$.
Introduction to NN and DL • An artificial neuron computes
 $y = f\left(\sum_{j=0}^m w_j x_j\right)$ with $x_0 = 1, w_0 = b$. Common activations: sigmoid $f(x) = \frac{1}{1+e^{-x}}$, ReLU $f(x) = \max(0, x)$.

• Binary output: 1 if $\mathbf{w}^T \mathbf{p} + b > 0$, else 0. Decision boundary: $\mathbf{w}^T \mathbf{p} + b = 0$. +1 from negative to positive.
• OR Gate Perceptron Design: Given $\mathbf{w} = [1, 1]^T$, point $\mathbf{p} = [0, 0.5]^T$ lies on the boundary. Then $1 \cdot 0 + 1 \cdot 0.5 + b = 0 \Rightarrow b = -0.5$.
• AND Gate Perceptron Design: Given $\mathbf{w} = [1, 1]^T$, point $\mathbf{p} = [0, 1.5]^T$ lies on the boundary. Then $1 \cdot 0 + 1 \cdot 1.5 + b = 0 \Rightarrow b = -1.5$.
• XOR and Multi-Layer Perceptron: Single-layer fails on XOR. Use hidden layer: Neuron 1: $\mathbf{w}_1 = [1, 1]^T, b = -0.5$; Neuron 2: $\mathbf{w}_2 = [-1, -1]^T, b = 1.5$. Output neuron AND: $\mathbf{w}_{\text{out}} = [1, 1], b = -1.5$.
• Forward Propagation: Given $W_1, W_2, W_3, b_1, b_2, b_3$, input \mathbf{p} , compute: $\mathbf{h}_1 = \sigma(W_1 \mathbf{p} + b_1), \mathbf{h}_2 = \sigma(W_2 \mathbf{h}_1 + b_2), \hat{y} = \sigma(W_3 \mathbf{h}_2 + b_3)$.
• Universal Approximation Theorem: (1) Approximating Arbitrary Decision Regions with 3-Layer MLP; Approximating a Function with 2-Layer MLP. (2) A hidden layer of nonconstant, bounded, and continuous neurons (with weights and biases), plus a linear output neuron (with weights but without bias) can approximate any continuous function to any arbitrary accuracy if there are enough hidden neurons.

Feedforward Computation
A multilayer NN maps input x to output y via: $y = f(x) = \sigma(W_L \sigma(W_{L-1} \dots \sigma(W_1 x + b_1) \dots + b_{L-1}) + b_L)$. All layers use weight matrices W_i and bias vectors b_i .
NN terminology Basics
• Classification: last layer = softmax, loss = cross-entropy.
• Regression: last layer = linear, loss = MSE.
• "2-layer NN(1-hidden-layer NN)" = input + 1 hidden + output.
• "3-layer NN(2-hidden-layer NN)" = input + 2 hidden + output.

Activation Functions
• Sigmoid: $\sigma(x) = \frac{1}{1+e^{-x}}$, output in $[0, 1]$, saturates.
• Tanh: $\tanh(x) = 2\sigma(2x) - 1$, output in $[-1, 1]$, zero-centered.
• ReLU: $\max(0, x)$, avoids vanishing gradient, fast to compute.
• Leaky ReLU: $f(x) = \alpha x$ if $x < 0$, else x , avoids dead neurons.
• Linear: $f(x) = cx$, used in regression output layer. Also known as identity activation function.
• Softmax: $\text{softmax}(z_i) = \frac{e^{z_i}}{\sum_j e^{z_j}}$, output is probability

distribution. • Maxout. • Exponential Linear Unit (ELU).
Data Preprocessing
• Standardization: zero mean, unit variance. $x' = \frac{x - \mu}{\sigma}$
• Normalization: scale to $[0, 1]$ or $[-1, 1]$: $x' = \frac{x - \min(x)}{\max(x) - \min(x)}$
Loss Functions
• Classification: Cross-entropy loss: $\mathcal{L}(\theta) = -\sum_{k=1}^K y_k \log \hat{y}_k$.
• Regression:
Mean squared error (MSE): $\mathcal{L}(\theta) = \frac{1}{n} \sum_i (y_i - \hat{y}_i)^2$.
Mean absolute error (MAE): $\frac{1}{n} \sum_i |y_i - \hat{y}_i|$

Gradient Descent (GD)
1. **Initialization**: Randomly initialize the model parameters θ^0 . Set $\theta^{\text{old}} = \theta^0$.
2. Compute the gradient of the loss function at θ^{old} : $\nabla \mathcal{L}(\theta^{\text{old}})$.
3. Update the parameters: $\theta^{\text{new}} = \theta^{\text{old}} - \alpha \nabla \mathcal{L}(\theta^{\text{old}})$ where α is the learning rate.
4. Set $\theta^{\text{old}} = \theta^{\text{new}}$ and return to step 2 (repeat until terminating).
GD does not guarantee reaching a global minimum.
• Backpropagation: Combines forward pass + backward pass using chain rule to compute gradients of loss w.r.t. weights.
• Batch GD: compute gradients over full dataset.
• Mini-batch GD: over subset (e.g., 32 to 256 samples). Mini-batch gradient approximates the full training set gradient well.
• SGD: mini-batch size = 1, fast but noisy. Less used.

GD with Momentum:
• Momentum: $v_t = \beta v_{t-1} + (1 - \beta) \nabla \mathcal{L}$, improves convergence.
• Nesterov Momentum: lookahead at next step before computing gradient. (Skip)
• Adam: computes a weighted average of past gradients and weighted average of past squared gradients. • Other: RMSprop, Adagrad, Adadelta, Nadam, etc. • Most used: Adam, SGD with momentum.

Learning Rate (LR)
• Too small \rightarrow slow convergence; Too large \rightarrow divergence or oscillations. • LR scheduling: reduce α over time. Approaches: step decay, exponential/cosine decay, reduce by constant, warmup/increase then cool down).

Underfitting and Overfitting
• Underfitting: high train + validation error (too simple).
• Overfitting: low train error, high validation error (too complex).
Regularization Techniques
• ℓ_2 regularization: $\mathcal{L}_{\text{reg}} = \mathcal{L}(\theta) + \lambda \sum_k \theta_k^2$, penalty for large.
• ℓ_1 regularization: $\sum_k |\theta_k|$, perform worse than ℓ_2 .
• Elastic net: mix of ℓ_1 and ℓ_2 . • Dropout: randomly deactivate neurons during training (e.g., $p = 0.5$). • Early stopping: monitor validation loss, stop if no improvement after n epochs.

Batch Normalization • Normalize each mini-batch: $\hat{x} = \frac{x - \mu}{\sigma}$
• Less dependent on initialization. Normalizes activations. Reduces vanishing/exploding gradients. Reduces internal covariate shift.

Stabilizes training, enables higher learning rates. Common in CNNs.
Hyperparameter Tuning • layer size, lr(schedule), optimizer, regularization, batch size, activation, loss f.
• Methods: grid search(check all with step), random search, Bayesian optimization. Use k-fold cross-validation if data is limited.
Why MLP Fails on Images MLPs are not translation invariant: shifting the input changes the output drastically. CNNs solve this with local receptive fields and shared weights.
Convolution Operation • 1D: Convolution flips and slides a kernel w over input x , $(x * w)[n] = \sum_{k=-\infty}^{\infty} x[k]w[n - k]$
 $x = [1, 2, 3, 4], w = [5, 6, 7] \rightarrow z = [5, 16, 34, 52, 45, 28]$
• 2D: Apply 2D filter over 2D image by elementwise multiplication and sum. Edge detectors highlight directional features.
Convolutional Neural Networks(CNNs) • Parameter sharing via convolution filters. • Sparse connections via local receptive fields. • Translation invariance. • Fewer parameters than MLPs: faster train.
Residual Networks (ResNets) • Use identity skip connections: output = layer(x) + x. • Help prevent vanishing gradients.
• Enable very deep models (e.g., 18, 50, 152 layers).
CNN Details • Each filter spans full input depth. • Produces one activation map per filter. • Multiple filters produce stacked output (e.g., 6 filters \rightarrow output depth = 6).
CNN Architecture Typical: $[(\text{CONV} - \text{ReLU}) \times N - \text{POOL}] \times M - [\text{FC} - \text{ReLU}] \times K - \text{Softmax}$. Modern trends: deep models (VGG, ResNet), small filters (3×3) , less pooling.
Spatial Dimensions For input of size $W_1 \times H_1 \times C$, filter size F , stride S , padding P , number of filters K : it will produce an output of $W_2 \times H_2 \times K$
 $W_2 = \frac{W_1 - F + 2P}{S} + 1$. Same for H_2 .
Number of parameters per filter: $F^2 C + 1$ (bias), total: $K(F^2 C + 1)$ for K filters.
Example Calculation: Input: $32 \times 32 \times 3$, Filters: 10 of size 5×5 , stride 1, padding 2 Output spatial size: $(32 + 2 * 2 - 5)/1 + 1 = 32 \rightarrow 32 \times 32 \times 10$ Params: $5 * 5 * 3 + 1 = 76$ per filter \rightarrow total 760 parameters.
Padding Strategy Zero-padding preserves spatial size. To preserve size: use padding $P = \frac{F-1}{2}$ if $S = 1$
Pooling Layer • Reduces spatial size, parameters, overfitting.
• Max pooling: keep largest value in window. • Avg pooling: keep mean value. • Output size: $W_2 = \frac{W_1 - F}{S} + 1$, same for H_2 .
• 0 learnable parameters.
Flatten and FC Layer • Flatten: convert activation map to 1D before FC. • FC: standard dense layer connects to all inputs.

```
import torch
X = torch.tensor([[0, 0], [0, 1], [1, 0], [1, 1]], dtype=torch.float)
Y = torch.tensor([0, 1, 1, 0], dtype=int)
network = torch.nn.Linear(2, 2)
loss_fn = torch.nn.CrossEntropyLoss()
optim = torch.optim.SGD(network.parameters(), lr=0.1)
epochs = 50
for epoch in range(epochs):
    for x, y in zip(X, Y):
        optim.zero_grad()
        y_pred = network(x)
        loss = loss_fn(y_pred, y)
        loss.backward()
        optimizer.step()

import numpy as np
class NeuralNetwork:
    def __init__(self, input_size, hidden_size, output_size, learning_rate, momentum=None):
        self.input_size = input_size
        self.hidden_size = hidden_size
        self.output_size = output_size
        self.learning_rate = learning_rate
        self.momentum = momentum
        self.weights_input_hidden = np.random.rand(input_size, hidden_size)
        self.bias_hidden = np.zeros((1, hidden_size))
        self.weights_hidden_output = np.random.rand(hidden_size, output_size)
        self.bias_output = np.zeros((1, output_size))
        if self.momentum:
            self.velocity_input_hidden = np.zeros((input_size, hidden_size))
            self.velocity_hidden_output = np.zeros((hidden_size, output_size))
    def sigmoid(self, x):
        return 1 / (1 + np.exp(-x))
    def sigmoid_derivative(self, x):
        return x * (1 - x)
    def forward_pass(self, X):
        self.hidden_output = self.sigmoid(np.dot(X, self.weights_input_hidden) + self.bias_hidden)
        self.output = self.sigmoid(np.dot(self.hidden_output, self.weights_hidden_output) + self.bias_output)
        return self.output
    def backward_pass(self, X, y, output):
        error = y - output
        output_delta = error * self.sigmoid_derivative(output)
        error_hidden = output_delta.dot(self.weights_hidden_output.T)
        hidden_delta = error_hidden * self.sigmoid_derivative(self.hidden_output)
        if self.momentum:
            self.velocity_input_hidden = self.momentum * self.velocity_input_hidden + X.T.dot(hidden_delta) * self.learning_rate
            self.velocity_hidden_output = self.momentum * self.velocity_hidden_output + self.hidden_output.T.dot(output_delta) * self.learning_rate
        self.weights_input_hidden += X.T.dot(hidden_delta) * self.learning_rate
        self.bias_output += np.sum(output_delta, axis=0, keepdims=True) * self.learning_rate
        self.bias_hidden += np.sum(hidden_delta, axis=0, keepdims=True) * self.learning_rate
    def train(self, X, y, epochs):
        for epoch in range(epochs):
            output = self.forward_pass(X)
            self.backward_pass(X, y, output)
            if epoch % 1000 == 0:
                loss = np.mean(np.square(y - output))
                print(f'Epoch {epoch}, Loss: {loss:.4f}')
```