Speech Signal Analysis

DEF (Dithering) Dithering adds noise to the signal which avoids the signal having zeros (which may be problematic when taking the log).

$$y[t] = x[t] + \epsilon[t]$$
 $\epsilon[t] \sim \mathcal{N}(0, \sigma^2)$

DEF (DC Offset) Most processing techniques assume that the signal is centered around 0.

$$y[t] = x[t] - \frac{1}{T}\mu_x$$
 $\mu_x = \frac{1}{T}\sum_{t=1}^{T}x[t]$

DEF (*Pre-emphasis*) Emphasize high frequency components.

$$y[t] = x[t] - \alpha x[t-1]$$

where α is a constant (usually 0.97)

DEF (Discrete Fourier Transform (DFT)) The dot product of the signal with the sinusoids of the frequency (Fourier basis).

$$X[k] = \sum_{t=0}^{T-1} x[t]e^{i2\pi kt/T}$$

DFT decomposes the signal into frequency components. Commonly ${\mathcal F}$ is used to denote the DFT operator.

THM (Properties of DFT)

- Linearity: $\mathcal{F}(a_1x_1 + a_2x_2) = a_1\mathcal{F}(x_1) + a_2\mathcal{F}(x_2)$ Shift: If y[t] = x[t-m], then $Y[k] = e^{i2\pi km/T}X[k]$

DEF (Windowing) Given a signal x[t], a window function w[t] (Hamming, Hann, etc.)

$$y[t] = x[t]w[t]$$

DEF (Short-Time Fourier Transform (STFT)) Given a signal x[t], a window function w[t], and a hop size h, the STFT is given by:

$$X[k,m] = \sum_{t=0}^{T-1} x[t]w[t-m]e^{-i2\pi kt/T}$$

where m is the current frame index and k is the frequency index. This gives a complex spectrogram of the signal.

DEF (Complex Spectrogram) Given a frequency bin X[k] = a + bi,

- Real: $\mathcal{R}(X[k]) = a$
- Imaginary: $\Im m(X[k]) = b$
- Magnitude: $|X[k]| = \sqrt{a^2 + b^2}$
- Phase: $\angle X[k] = \arccos \frac{a}{\sqrt{a^2+b^2}}$ Energy: $P[k] = |X[k]|^2$

DEF (Mel Spectrogram) After obtaining the spectrogram, X[k], and using the mel filter bank, $H_n[k]$, the mel spectrogram is given by:

$$Y[k, n] = \sum_{k=0}^{K-1} |X[k]| H_n[k]$$

DEF (Mel Frequency Cepstral Coefficients (MFCCs)) A common feature extraction technique for speech recognition which is able to capture speaker characteristics and phonetic content.

- 1. Extract the mel spectrogram Y[k, n]
- 2. Apply the Discrete Cosine Transform (DCT) to the mel spectrogram.
- 3. Take the first p coefficients (usually 13).

Hidden Markov Models

DEF (Objective of ASR) Given a sequence of acoustic feature vectors X, and W denotes a word sequence, ASR aims to find the most likely word sequence W^*

$$W^* = \operatorname*{arg\,max}_{W} P(W|X)$$

COR (Decomposition of P(W|X)) We can decompose P(W|X) with Bayes' theorem:

$$P(W|X) = \frac{P(X|W)P(W)}{P(X)} \propto P(X|W)P(W)$$

P(X|W) is the **Acoustic Model** and P(W) is the **Language** Model.

DEF (Modelling the Acoustic Model with HMMs) Commonly, the left-to-right HMM is used to model the acoustic model. As a word is composed of multiple phonemes, we can model each phoneme with a left-to-right HMM and then concatenate them to form a HMM for the word. For the phoneme HMMs, we typically use a left-to-right HMM with 3 states which as a consequence also enforces a minimum phone duration.

DEF (The three fundamental problems of HMMs)

1. **Likelihood** - Determine the overall likelihood of an observation sequence $X = x_1, x_2, \dots, x_T$ given an HMM topology \mathcal{M} :

$$P(X|\mathcal{M})$$

- 2. **Decoding** Determine the most likely sequence of hidden states $Q = q_1, q_2, \dots, q_T$ given an observation sequence X = x_1, x_2, \ldots, x_T and an HMM topology \mathcal{M}
- 3. **Training** Given an observation sequence $X = x_1, x_2, \dots, x_T$ and an HMM topology \mathcal{M} , determine the optimal state occupation probabilities.

ALG (Forward Probability $[\alpha_i(t)]$) The **Likelihood** and **Training** problem rely on forward probabilities. The forward probability $\alpha_i(t)$ is the probability of observing the observation sequence x_1, \dots, x_t and being in state j at time t. This can be computed recursively:

• Initialisation:

$$\alpha_j(0) = 1$$
 $j = 0$
 $\alpha_j(0) = 0$ $j \neq 0$

• Recursion:

$$\alpha_j(t) = \left(\sum_{i=0}^{J} \alpha_i(t-1)a_{ij}\right)b_j(x_t)$$

• Termination:

$$P(X|\mathcal{M}) = \alpha_E = \sum_{i=1}^{J} \alpha_i(T) a_{iE}$$

ALG (Viterbi Algorithm) The Decoding problem is solved by the Viterbi Algorithm. Define likelihood of the most probable partial path in state j at time t as $V_i(t)$ and the Viterbi backpointer as $B_i(t)$ (the most probable state at time t-1 that leads to state j at time t). Then the Viterbi algorithm can be described as:

• Initialisation:

$$V_0(0) = 1$$

 $V_j(0) = 0$ $j \neq 0$
 $B_j(0) = 0$

• Recursion:

$$V_j(t) = \max_{i=0}^{J} V_i(t-1)a_{ij}b_j(x_t)$$
$$B_j(t) = \arg\max_{i=0}^{J} V_i(t-1)a_{ij}$$

• Termination:

$$V_E = \max_{j=1}^{J} V_j(T) a_{jE}$$
$$B_E = \arg \max_{j=1}^{J} V_j(T) a_{jE}$$

ALG (Training: Hard Assignment with Viterbi Training) If we know the state-time alignment of the training data (i.e. the sequence of states that generated the observation sequence), then the transition and observation probabilities can be estimated as follows:

$$a_{ij} = \frac{C(i \to j)}{\sum_{k} C(i \to k)}$$
$$\mu_{j} = \frac{\sum_{t} z_{jt} x_{t}}{\sum_{t} z_{jt}}$$
$$b_{j} = \mathcal{N}(x_{t} | \mu_{j}, \Sigma_{j})$$

where z_{jt} is an indicator variable that is 1 if the t-th observation was generated by the j-th state. $C(i \to j)$ is the number of times the transition from state i to state j is made.

DEF (Backward Probability $[\beta_j(t)]$) The backward probability $\beta_j(t) = P(x_{t+1}, \dots, x_T | q_t = j, \mathcal{M})$ is the probability of observing the observation sequence from time t+1 to the end, given that the HMM is in state j at time t.

• Initialisation:

$$\beta_j(T) = a_{jE}$$

• Recursion:

$$\beta_i(t) = \sum_{j=1}^{J} a_{ij} b_j(x_{t+1}) \beta_j(t+1)$$

• Termination:

$$P(X|\mathcal{M}) = \beta_0(0) = \sum_{j=1}^{J} a_{0j}b_j(x_1)\beta_j(1) = \alpha_E$$

DEF (State Occupation Probability $[\gamma_j(t)]$) $\gamma_j(t)$ is the probability of being in state j at time t, given the complete observation sequence X and the model \mathcal{M} . It combines forward and backward probabilities.

$$\begin{split} \gamma_{j}(t) &= P(q_{t} = j | X, \mathcal{M}) \\ &= \frac{P(q_{t} = j, X_{1:T} | \mathcal{M})}{P(X | \mathcal{M})} \\ &= \frac{P(q_{t} = j, X_{1:t}, X_{t+1:T} | \mathcal{M})}{\alpha_{E}} \\ &= \frac{P(X_{t+1:T} | q_{t} = j, X_{1:t}, \mathcal{M}) \times P(q_{t} = j, X_{1:t} | \mathcal{M})}{\alpha_{E}} \\ &= \frac{P(X_{t+1:T} | q_{t} = j, \mathcal{M}) \times P(q_{t} = j, X_{1:t} | \mathcal{M})}{\alpha_{E}} \\ &= \frac{\beta_{j}(t) \alpha_{j}(t)}{\alpha_{E}} \end{split}$$

This represents the expected proportion of time spent in state j at time t. Note that $\sum_{j=1}^J \gamma_j(t) = 1$ for any t.

DEF (Transition Occupation Probability $(\xi_{ij}(t))$) $\xi_{ij}(t)$ is the probability of being in state i at time t and transitioning to state j at time t+1, given the complete observation sequence X and the model \mathcal{M} .

$$\begin{split} \xi_{ij}(t) &= P(q_t = i, q_{t+1} = j | X, \mathcal{M}) \\ &= \frac{P(q_t = i, q_{t+1} = j, X_{1:T} | \mathcal{M})}{P(X | \mathcal{M})} \\ &= \frac{P(X_{1:t}, q_t = i, q_{t+1} = j, X_{t+1:T} | \mathcal{M})}{\alpha_E} \\ &= \frac{\alpha_i(t) a_{ij} b_j(x_{t+1}) \beta_j(t+1)}{\alpha_E} \end{split}$$

ALG (*EM Algorithm (Baum-Welch) for HMM Training*) The Baum-Welch algorithm is an instance of the Expectation-Maximization (EM) algorithm used to find the maximum likelihood estimate of the parameters ($\lambda = \{A, B, \pi\}$) of an HMM when the state sequence is unknown. It iteratively performs two steps:

- E-step (Expectation): Using the current parameter estimates λ^{old} , compute the expected values needed for the M-step. This involves:
 - Computing forward probabilities $\alpha_i(t) \ \forall j, t$.
 - Computing backward probabilities $\beta_i(t) \forall j, t$.
 - Computing state occupation probabilities $\gamma_i(t) \ \forall j, t$
 - Computing transition occupation probabilities $\xi_{ij}(t) \forall i, j, t$
 - If using GMMs for b_j , compute component-state occupation probabilities $\gamma_{jm}(t)$ (see below).
- M-step (Maximization): Update the parameters λ^{new} to maximize the expected complete-data log-likelihood, using the expectations computed in the E-step. See specific update formulas below.
- Repeat E-step and M-step until the parameters or the likelihood $P(X|\lambda)$ converge.

DEF (*M-Step: Transition Probabilities* (a_{ij})) The transition probability a_{ij} is updated as the expected number of transitions from state i to j, divided by the expected total number of transitions *out* of state i.

$$\begin{split} \hat{a}_{ij}^{\text{new}} &= \frac{\text{Expected transitions } i \rightarrow j}{\text{Expected transitions out of } i} \\ &= \frac{\sum_{t=1}^{T-1} \xi_{ij}(t)}{\sum_{k=1}^{J} \sum_{t=1}^{T-1} \xi_{ik}(t)} = \frac{\sum_{t=1}^{T-1} \xi_{ij}(t)}{\sum_{t=1}^{T-1} \gamma_i(t)} \end{split}$$

For a corpus of R utterances, sum the numerators and denominators over all utterances $r=1,\ldots,R$.

DEF (M-Step: Observation Probabilities $(b_j(\mathbf{x}))$ [Single Gaussian]) If the observation probability for state j is modeled by a single multivariate Gaussian $b_j(\mathbf{x}) = \mathcal{N}(\mathbf{x}; \boldsymbol{\mu}_j, \boldsymbol{\Sigma}_j)$, the parameters are updated using weighted maximum likelihood, where the weights are the state occupation probabilities $\gamma_j(t)$.

$$\hat{\boldsymbol{\mu}}_{j}^{\text{new}} = \frac{\sum_{t=1}^{T} \gamma_{j}(t) \mathbf{x}_{t}}{\sum_{t=1}^{T} \gamma_{j}(t)}$$

$$\hat{\boldsymbol{\Sigma}}_{j}^{\text{new}} = \frac{\sum_{t=1}^{T} \gamma_{j}(t) (\mathbf{x}_{t} - \hat{\boldsymbol{\mu}}_{j}^{\text{new}}) (\mathbf{x}_{t} - \hat{\boldsymbol{\mu}}_{j}^{\text{new}})^{T}}{\sum_{t=1}^{T} \gamma_{j}(t)}$$

DEF (*M-Step: Observation Probabilities* $(b_j(\mathbf{x}))$ [*GMM*]) If $b_j(\mathbf{x})$ is an M-component GMM, $b_j(\mathbf{x}) = \sum_{m=1}^M c_{jm} \mathcal{N}(\mathbf{x}; \boldsymbol{\mu}_{jm}, \boldsymbol{\Sigma}_{jm})$, we need the component-state occupation probability $\gamma_{jm}(t)$.

• Component-State Occupation Probability $(\gamma_{jm}(t))$: This is the probability of being in state j and using component m at time t, given X. It's calculated in the E-step after $\gamma_j(t)$:

$$\gamma_{jm}(t) = P(q_t = j, c_t = m | X, \mathcal{M}^{\text{old}})$$

$$= \gamma_j(t) \times \frac{c_{jm}^{\text{old}} \mathcal{N}(\mathbf{x}_t; \boldsymbol{\mu}_{jm}^{\text{old}}, \boldsymbol{\Sigma}_{jm}^{\text{old}})}{b_j^{\text{old}}(\mathbf{x}_t)}$$

$$= \frac{\alpha_j(t)\beta_j(t)}{\alpha_E} \times \frac{c_{jm}^{\text{old}} \mathcal{N}(\mathbf{x}_t; \boldsymbol{\mu}_{jm}^{\text{old}}, \boldsymbol{\Sigma}_{jm}^{\text{old}})}{\sum_{k=1}^{M} c_{jk}^{\text{old}} \mathcal{N}(\mathbf{x}_t; \boldsymbol{\mu}_{jk}^{\text{old}}, \boldsymbol{\Sigma}_{jk}^{\text{old}})}$$

• GMM Parameter Updates (M-step): Use $\gamma_{jm}(t)$ as weights for component m of state j.

$$\hat{c}_{jm}^{\text{new}} = \frac{\sum_{t=1}^{T} \gamma_{jm}(t)}{\sum_{k=1}^{M} \sum_{t=1}^{T} \gamma_{jk}(t)} = \frac{\sum_{t=1}^{T} \gamma_{jm}(t)}{\sum_{t=1}^{T} \gamma_{j}(t)}$$

$$\hat{\boldsymbol{\mu}}_{jm}^{\text{new}} = \frac{\sum_{t=1}^{T} \gamma_{jm}(t) \mathbf{x}_{t}}{\sum_{t=1}^{T} \gamma_{jm}(t)}$$

$$\hat{\boldsymbol{\Sigma}}_{jm}^{\text{new}} = \frac{\sum_{t=1}^{T} \gamma_{jm}(t) (\mathbf{x}_{t} - \hat{\boldsymbol{\mu}}_{jm}^{\text{new}}) (\mathbf{x}_{t} - \hat{\boldsymbol{\mu}}_{jm}^{\text{new}})^{T}}{\sum_{t=1}^{T} \gamma_{jm}(t)}$$