MUSIC ANALYSIS WITH A BAYESIAN MODEL

Lilia Ben Baccar MS-DS, ENSAE Paris lilia.benbaccar@ensae.fr Erwan Rahis MS-DS, ENSAE Paris erwan.rahis@ensae.fr

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Abstract

The goal of this project is to use Bayesian inference to analyze a music piece. We consider a music piece as a sequential data and segment the piece using a non-parametric Bayesian model. The music, to be represented as a discrete sequence of observations is processed and then modeled using a Hidden Markov Chain model. The segmentation is inferred by the model and compared to our music-theoretic analysis. Our method is mainly a combination of the papers by Lu Ren and al. [2] and Emily B. Fox and al. [1].

1 Bayesian model

To solve this problem, the approach taken is a non-parametric Bayesian model based on the *Hidden Markov Model* and a prior generated by a *Hierarchical Dirichlet Process* (HDP-HMM). We want to segment a waveform into a set of time intervals with no *a priori* on the number of intervals. The sequence is modelled using the Hidden Markov Model with the states to categorise the interval. In other words, we want the piece to be divided into intervals and each interval will be labelled to a state that we will later on analyse. For example a music piece can be divided into parts that are: intro, chorus, pre-chorus, outro. Those parts can be easily identified by the human ear and brain (music-theoretic analysis). We will compare our analysis to the results of the model.

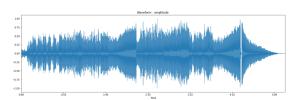


Figure 1: Audio waveform of the piece considered

1.1 Hidden Markov Model

The Hidden Markov Model used is described as follows: A discrete sequence of observed data $y = \{y_t\}_{t=1}^T$ with $y_t \in \{1, \dots, M\}$ and a corresponding hidden state sequence $s = \{s_t\}_{t=1}^T$ with $s_t \in \{1, \dots, I\}$. The states will be analysed as in Section 1 (intro, chorus, etc.). Each observation is associated to a state. The discrete HMM model is represented by its parameters $\theta = \{\pi, B, A\}$ such that:

- π is a set of state transition probabilities, the probability to transition from one state to another. For two states ρ and ξ , the transition probability from ρ to ξ is $\pi_{\rho\xi} = \mathbb{P}(s_{t+1} = \xi | s_t = \rho)$
- B is a set of emission probabilities, the probability that an observation associated to a state ρ is equal to $m: b_{\rho m} = \mathbb{P}(x_t = m | s_t = \rho)$
- A is the set of initial state distribution, the probability that the first observation in time t=1 is equal to the state ρ : $a_{\rho} = \mathbb{P}(s_1 = \rho)$

Because we will use a Hierchical Dirichlet Process to model the parameters θ , we will divide the sequence into J subsequences. This means that we will model J Hidden Markov Models. By taking the notation at the beginning of this section, we will have $j \in \{1, \ldots, J\}$ and $y = \{y_j\}_{j=1}^J$ with $y_j = \{y_{ji}\}_{i=1}^{N_j}$. Note that every subsequence x_j has a length N_j . We allow the parameters to be different in each subsequence such that:

$$x_j \sim HMM(\theta_j)$$

The joint distribution given θ is :

$$p(y|\theta) = \prod_{j=1}^J \left\{ \sum_{s_j} a_{s_{j,1}} \prod_{i=1}^{N_j-1} \pi_{s_{j,i+1}} \prod_{i=1}^{N_j} \ b_{s_{j,i}y_{j,i}} \right\}$$

1.2 Hierchical Dirichlet Process

A Hierchical Dirichlet Process is defined as a collection of Dirichlet Processes. Our collection of DP processes is the J parameters θ_j linked to the HMM. This kind of setting allows for analysing how inter-related are the subsequences one to the others. The first process of the collection is denoted as:

$$G_0 \sim DP(\gamma, H)$$

with γ a concentration parameter for the Dirichlet distribution which is symmetrical and H a base probability measure on Θ . The formal definition of a Dirichlet process is $\forall (B_k)_{k=1,...,K}$ finite partition of Θ :

$$(G_0(B_k))_{k=1,...,K}|\gamma, H \sim Dir((\gamma H(B_k))_{k=1,...,K})$$

We approach the Dirichlet Process with a stick-breaking process and we can write G_0 as follows:

$$G_0 = \sum_{k=1}^{\infty} \beta_k \delta_{\theta_k}$$

with δ_{θ_k} an indicator function in θ_k and β_k the mixture weights computed such that :

$$\beta_k \sim Stick(\gamma)$$

$$\beta_k = v_k \prod_{l=1}^{k-1} (1 - v_l)$$

$$v_k \sim Beta(1, \gamma)$$

with v_k the independent random variables following a Beta distribution. The process to obtain the probabilities β_k is pictured as breaking a unit-length stick because we start off with a stick of length 1 and at each step we we break a portion of the what is remaining stick according to v_k and assign to β_k the piece broken. The smaller the α the more concentrated is the distribution. The first process G_0 (first subsequence) is then used as the base measure for the other processes (other subsequences):

$$G_{j} \sim DP(\alpha, G_{O})$$

$$G_{j} = \sum_{t=1}^{\infty} \tilde{\pi}_{jt} \delta_{\theta_{jt}^{\star}}$$

$$\tilde{\pi}_{j} | \alpha \sim Stick(\alpha)$$

$$\theta_{jt}^{\star} | G_{0} \sim G_{0}$$

$$\theta_{ji}^{\prime} | G_{j} \sim G_{j}$$

$$y_{ji} | \theta_{ji}^{\prime} \sim F(\theta_{ij}^{\prime})$$

with $j \in \{1, ..., J\}$, $i \in \{1, ..., N_j\}$. We know that multiple θ_{jt}^* can take the same value θ_k . We can write G_j as a function of unique states:

$$G_{j} \sim \sum_{k=1}^{\infty} \pi_{jk} \delta_{\theta_{k}}$$
$$\pi_{j} | \alpha, \beta \sim DP(\alpha, \beta)$$
$$\theta_{k} | H \sim H$$

with π_{jk} the distribution of the subsequence j over the state k with :

$$\pi_{jk} = \sum_{t \mid \theta_{jt}^{\star} = \theta_k} \tilde{\pi}_{jt}$$

The model can be re written with an indicator variable z_{ii} for the state assigned to the observation y_{ii} :

$$G_{j} \sim DP(\alpha, G_{O})$$

$$z_{ji}|\pi_{j} \sim \pi_{j}$$

$$y_{ji}|\{\theta_{k}\}, z_{ij} \sim F(\theta_{z_{ii}})$$

With $F(\cdot)$ the fixed distribution of the data sample. According to [1], we compute a sticky version of the model by adding a κ parameter that will slow down the transitions between states and avoid having unrealistically fast dynamics as results. The new transition distributions are modelled as follows:

$$\beta | \gamma \sim Stick(\gamma)$$

$$\pi_j | \alpha, \kappa, \beta \sim DP\left(\alpha + \kappa, \frac{\alpha\beta + \kappa\delta_j}{\alpha + \kappa}\right)$$

By adding this parameter, we are increasing the probability of self-transition by an amount proportional to k. If the probability that an observation has the same state as the one after is higher, the changing dynamics will slow down. That is why the κ parameter is added on the j^{th} component of $\alpha\beta$.

2 Computational method

We use the Gibbs Block Sampler [1] for the inference of the parameters of the sticky HDP-HMM. The first step is to sample β and π . The sampler is defined with a finite Dirichlet prior on θ_j such that :

$$\beta | \gamma \sim Dir(\gamma/L, \dots, \gamma/L)$$

 $\pi_j | \alpha, \beta \sim Dir(\alpha\beta_1, \dots, \alpha\beta_L)$

This finite approximation is made because when $L \to \infty$ then it converges to the HDP mixture model. The derived posterior distribution is :

$$\beta|\bar{n}, \gamma \sim Dir(\gamma/L + \bar{m}_{\cdot 1}, \dots, \gamma/L + \bar{m}_{\cdot L})$$

$$\pi_{j}|z_{1:T}, \alpha, \beta \sim Dir(\alpha\beta + n_{j1}, \dots, \alpha\beta_{j} + \kappa + n_{jj}, \dots, \alpha\beta_{L} + n_{jL})$$

with n_{jk} the number of j to k transitions in the state sequence j. The choice of initialisation of the base measure is made regarding the paper with :

$$y|z = k \sim \mathcal{N}(\mu_k, \sigma_k^2)$$

$$\mu_k \sim \mathcal{N}(0, s^2)$$

$$\sigma_k^2 \sim InvGamma(a, b)$$
(1)

After, choosing the hyper-parameters and base measure, the next step is to sample the hidden states z_t . This is done with a forward-backward algorithm to jointly sample z given y for the sticky HDP-HMM. The conditional distribution is :

$$p(z|y, \pi, \theta) = p(z_t|z_{t-1}, y, \pi, \theta)p(z_{t-1}|z_{t-2}, y, \pi, \theta) \dots p(z_2|z_1y, \pi, \theta)p(z_1|y, \pi, \theta)$$

The algorithm starts by sampling z_1 from $p(z_1|y,\pi,\theta)$ then sample z_2 from $p(z_2|z_1,y,\pi,\theta)$ and so on. The conditional distribution is :

$$p(z_1|y, \pi, \theta) \propto p(z_1) f(y_1|\theta_{z_1}) \sum_{z_2, T} \prod_t p(z_t|\pi_{z_{t-1}}) f(y_t|\theta_{z_t})$$

$$\propto p(z_1) f(y_1|\theta_{z_1}) \sum_{z_2} p(z_2|\pi_{z_1}) f(y_2|\theta_{z_2}) m_{3,2}(z_2)$$

$$\propto p(z_1) f(y_1|\theta_{z_1}) m_{2,1}(z_1)$$

where $m_{t,t-1}(z_{t-1})$ is called the bakward message passed from z_t to z_{t-1} and defined by:

$$m_{t,t-1}(z_{t-1}) \propto p(y_{t:T}|z_{t-1},\pi,\theta)$$

In the algoritm the messages are first initialized to :

$$m_{T+1,T}(k) = 1$$

and then with a backward loop for each $k \in \{1, ..., L\}$ we compute the messages :

$$m_{t,t-1}(k) = \sum_{j=1}^{L} \pi_k(j) \mathcal{N}(y_t, \mu_j, \Sigma_j) m_{t+1,t}(j)$$
$$= \sum_{j=1}^{L} \pi_k(j) \frac{1}{\Sigma_j} \exp\left[-\frac{1}{2} \left(\frac{y_t - \mu_j}{\Sigma_j}\right)^2\right] m_{t+1,t}(j)$$

Because of (1) we have $\mathcal{N}(y_t, \mu_j, \Sigma_j) \propto \frac{1}{\Sigma_j} e^{-\frac{1}{2} \left(\frac{y_t - \mu_j}{\Sigma_j}\right)^2}$. Then we derive the general distribution for z_t which is:

$$p(z_t|z_{t-1}, y, \pi, \theta) \propto p(z_t|\pi_{z_{t-1}}) f(y_t|\theta_{z_t}) m_{t+1,t}(z_t)$$

The algorithm computes the z_t with a forward loop starting with $n_{jk} = 0$ and the cluster $\mathcal{Y}_k = \emptyset$ and for each $(j,k) \in \{1,\ldots,L\}^2$ computes the probability of y_t :

$$f_k(y_t) = \pi_{z_{t-1}}(k) \mathcal{N}(y_t, \mu_k, \Sigma_k) m_{t+1,t}(j)$$

and then assigns a state z_t with :

$$z_t \sim \sum_{k=1}^{L} f_k(y_t) \delta(z_t, k)$$

with $\delta(z_t, k) = \begin{cases} 1 & z_t = k \\ 0 & z_t \neq k \end{cases}$ and eventually updates n_{z_t-1, z_t} and adds y_t to the cluster \mathcal{Y}_k for a new assignment $z_t = k$.

3 Data processing

We consider an acoustic signal in a mp3 format loaded in Python using the librosa package. The sampling rate is 22.05kHz which means we have 22050 observations par second of audio signal. The acoustic signal is first loaded as a waveform or amplitude as shown in Figure 1. It represents the pressure change recorded by a recording instrument such as a microphone. We then divide the piece into 100ms contiguous frames. Each amplitude frame is processed by computing its Mel frequency sceptral coefficients (MFCCs). These coefficients are derived from the Melspectrogram that is a log-scale of pitches judged to be equally distanced to one another. The Mel spectrogram is the representation showing the mel frequencies according to time and intensity. The MFCCs are the most popular feature used in audio processing. We compute 40 columns of coefficients for each frame. To obtain a vector sequential data, we quantize each frame of coefficient. First, we whiten the columns so each one of them have a unit variance. Then we use the k-means algorithm with 16 codebooks and pass them to the vector quantization algorithm. The result is a 1D discrete sequential data representing the coefficients.

4 Results

References

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