Bayesian Stochastic Partition Models For Markovian Dependence Structures

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Academic dissertation

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<table>
<thead>
<tr>
<th>Role</th>
<th>Name</th>
<th>Institution</th>
<th>Location</th>
</tr>
</thead>
<tbody>
<tr>
<td>Supervisor</td>
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<td>Department of Mathematics and Statistics</td>
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<td></td>
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<td></td>
</tr>
<tr>
<td>Pre-examiners</td>
<td>Assistant Professor Harri Lähdesmäki</td>
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<td>Finland</td>
</tr>
<tr>
<td></td>
<td></td>
<td>Aalto University</td>
<td></td>
</tr>
<tr>
<td></td>
<td>Professor Jaakko Nevalainen</td>
<td>School of Health Sciences</td>
<td>Finland</td>
</tr>
<tr>
<td></td>
<td></td>
<td>University of Tampere</td>
<td></td>
</tr>
<tr>
<td>Custos</td>
<td>Professor Jukka Corander</td>
<td>Department of Mathematics and Statistics</td>
<td>Finland</td>
</tr>
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<td></td>
<td></td>
<td>University of Helsinki</td>
<td></td>
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<tr>
<td>Opponent</td>
<td>Reader Korbinian Strimmer</td>
<td>Department of Epidemiology and Biostatistics</td>
<td></td>
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<td></td>
<td>Imperial College London</td>
<td>United Kingdom</td>
</tr>
</tbody>
</table>
Abstract

In various fields of knowledge we can observe that the availability of potentially useful data is increasing fast. A prime example is the DNA sequence data. This increase is both an opportunity and a challenge as new methods are needed to benefit from the big data sets. This has sparked a fruitful line of research in statistics and computer science that can be called machine learning. In this thesis, we develop machine learning methods based on the Bayesian approach to statistics. We address a fairly general problem called clustering, i.e. dividing a set of objects to non-overlapping group based on their similarity, and apply it to models with Markovian dependence structures. We consider sequence data in a finite alphabet and present a model class called the Sparse Markov chain (SMC). It is a special case of a Markov chain (MC) model and offers a parsimonious description of the data generating mechanism. A Variable length Markov chain (VLMC) is a popular sparse model presented earlier in the literature and it has a representation as an SMC model. We develop Bayesian clustering methodology for learning the SMC and other Markovian models.

Another problem that we study in this thesis is causal inference. We present a model and an algorithm for learning causal mechanisms from data. The model can be considered as a stochastic extension of the sufficient-component cause model that is popular in epidemiology. In our model there are several causal mechanisms each with its own parameters. A mixture distribution gives a probability that an outcome variable is associated with a mechanism.

Applications that are considered in this thesis come mainly from computational biology. We cluster states of Markovian models estimated from DNA sequences. This gives an efficient description of the sequence data when comparing to methods reported in the literature. We also cluster DNA sequences with Markov chains, which results in a method that can be used for example in the estimation of bacterial community composition in a sample from which DNA is extracted. The causal model and the related learning algorithm are able to estimate mechanisms from fairly challenging data. We have developed the learning algorithms with big data sets in mind. Still, there is a need to develop them further to handle ever larger data sets.
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Helsinki, September 2014

Väinö Jääskinen
List of original articles

This thesis consists of four articles and an introductory part. We refer to the articles by Roman numerals I-IV.


Author’s contribution to Articles I-IV

I. VJ contributed to developing the method and had the main responsibility for the implementation and empirical testing, while JX also contributed to these. VJ also participated in writing the article, while JC had the main role.

II. VJ and VP contributed equally to the implementation, empirical results and writing the article. The method was mainly developed by JC and VP.

III. The method was jointly developed by the authors. VJ participated in the implementation, empirical testing and writing the article, while JX had the main responsibility for these.

IV. VJ participated in developing the method and writing the paper, while TK had the main role. VJ contributed to the implementation and testing the model, while HW had the main responsibility for these and empirical testing. This paper was part of the PhD thesis of Helga Westerlind at Karolinska Institut.
## Contents

Acknowledgements

<table>
<thead>
<tr>
<th>Chapter</th>
<th>Title</th>
<th>Page</th>
</tr>
</thead>
<tbody>
<tr>
<td>1</td>
<td>Introduction</td>
<td>1</td>
</tr>
<tr>
<td>2</td>
<td>Bayesian Statistics</td>
<td>5</td>
</tr>
<tr>
<td>2.1</td>
<td>The Bayesian Approach</td>
<td>5</td>
</tr>
<tr>
<td>2.2</td>
<td>Modeling</td>
<td>6</td>
</tr>
<tr>
<td>3</td>
<td>Markovian Models</td>
<td>9</td>
</tr>
<tr>
<td>3.1</td>
<td>Markov Chain</td>
<td>9</td>
</tr>
<tr>
<td>3.2</td>
<td>Variable Length Markov Chain</td>
<td>11</td>
</tr>
<tr>
<td>3.3</td>
<td>Mixture Transition Distribution</td>
<td>12</td>
</tr>
<tr>
<td>3.4</td>
<td>Sparse Markov Chain</td>
<td>13</td>
</tr>
<tr>
<td>4</td>
<td>Graphical Representations</td>
<td>15</td>
</tr>
<tr>
<td>5</td>
<td>Bayesian Clustering</td>
<td>19</td>
</tr>
<tr>
<td>5.1</td>
<td>Prior and Likelihood</td>
<td>19</td>
</tr>
<tr>
<td>5.2</td>
<td>Expectation-Maximization Algorithm</td>
<td>24</td>
</tr>
<tr>
<td>5.3</td>
<td>Classification EM Algorithm</td>
<td>25</td>
</tr>
<tr>
<td>5.4</td>
<td>Stochastic Search</td>
<td>27</td>
</tr>
<tr>
<td>5.5</td>
<td>Recursive Search</td>
<td>28</td>
</tr>
</tbody>
</table>
6 Causal Inference 31
  6.1 Sufficient-Component Cause Model 32
  6.2 Do-Calculus 34

7 Discussion 39

Bibliography 43
Introduction

One characteristic of the contemporary world is the abundance of potentially relevant data in various fields. This holds for scientific research, commercial applications and technology in general. There is a pressing need for new methods in data analysis. In statistics, prevailing attitude to big data in all of its forms has been somewhat ambivalent. This is because a given statistical method usually works for data sets that are large enough but not too large. In this thesis, the theme of large data sets and how to handle them is explored. Statistics is not the only discipline that is concerned with large data sets, as computer science is also intimately involved in these matters. Indeed, the relatively new field of machine learning combines perspectives from both statistics and computer science.

In this thesis, we start with the core ideas of Bayesian statistics and apply them to machine learning in applications including computational biology, especially analysis of genome sequence variation. This is natural as the Bayesian approach to statistics has become important for machine learning (Bishop, 2006). Bayesian statistics has a long and interesting history starting from Thomas Bayes in the 18th century (Bernardo and Smith, 1994).

A central motivation for this work comes from sequential data as it enables the use of Markovian models in their rich variety. These models have a history of over hundred years. The first definition and also application of Markov chains was by Andrei Markov to model probabilities of vowels and consonants in Alexander Pushkin’s verse novel *Eugene Onegin* (Hayes,
2013). As an example of sequence data we have in Article II 91 240 DNA sequences each having after preprocessing an approximate length of 500 base pairs. Another issue we consider is the study of causal inference, especially in the context of epidemiology. For example, we can think of a situation where a disease has two causing mechanisms and the effective mechanism is chosen randomly for each individual. Then the probability of catching the disease depends on the chosen mechanism and the covariates, for example age and sex. One unifying theme in this thesis is the use of classification, clustering and probabilistic reasoning to make sense of large data sets.

Here, we present some facts about DNA as it is central for applications described in this thesis. In DNA (deoxyribonucleic acid), there are four types of nucleobases: adenine (A), guanine (G), thymine (T) and cytosine (C) (Kimball, 1994). These are nitrogen-containing ring-like structures. Bases take part in forming polymers of nucleotides. DNA itself is such a polymer of nucleotides and it has two strands giving it the double helix structure. The genetic code of the DNA includes information that is used in the synthesis of proteins and thus it controls development of living organisms together with environmental factors.

The following chapters provide background to the four articles and highlight important issues. In Chapter 2, the Bayesian approach to statistics is presented. Both key equations and modeling principles are covered. In Chapter 3, Markovian models are described. There are several models that each share a form of Markov property. Chapter 4 includes graphical representations of the Markovian models. These illustrate the similarity and differences between Markovian models. In Chapter 5, the framework of Bayesian clustering is introduced. This chapter includes a variety of concepts that are needed for modeling as well as description of algorithms. Chapter 6 describes causal inference from both epidemiological and probabilistic point of view, showing examples of theories related to causality. Finally, Chapter 7 presents conclusions based on the articles and other material as well as points directions for future research.

Here, we describe some of the notation used in this thesis. Overall, notation should be clear from the context. Notation in each chapter is chosen to bring out the subject matter with clarity. This leads to a situation where
there are several symbols for concepts that are similar and on the other hand multiple uses for a single symbol. This clarification is provided to help the reader.

Firstly, there are several symbols for data. In Chapter 2, $X$ can be considered as a random variable. Its dimension is not defined explicitly. In Chapter 3, $X_n$ is a random variable and $n$ is an index variable. A realization of $X_n$ is denoted by $j_n$ or $x_n$. In Chapter 5, $x$ denotes a set of observations with $n$ data objects. The dimension of the data objects varies depending on the context. $x$ represents the observed data in the context of expectation-maximization (EM) algorithm.

Generally, $\theta$ represents collectively quantitative parameters. Also its dimension varies depending on the context. In the context of Markov chains, $n$ means the total number of observed transitions. In the context of clustering, it can also mean the number of data objects to be clustered. $S$ denotes a partition variable. In Chapter 2, both $H$ both and $A$ are taken to mean a hypothesis.

In the description of the classification EM algorithm, $M$ is the number of mechanisms and $N$ the number of observations. $j$ denotes an index of a mechanism and $l$ an index of an observation (subject). $y^{(l)}$ denotes the health status of an observed subject $l$ (ill or not) and $x^{(l)}$ denotes the covariate vector of the same subject. Assignment of each observation to a mechanism is denoted by $u^{(l)}$. An alternative representation is given by indicator variables of the form $u_j^{(l)}$. 
2.1 The Bayesian Approach

Here, we present some central aspects of Bayesian statistics as a background for later chapters. Generally, in statistical inference we are concerned with hypothesis about quantitative parameters \( \theta \) and we want to assess these hypothesis based on empirical data \( X \). This means using data to draw conclusions about unobserved quantities (Gelman et al., 2004).

A fundamental result for Bayesian statistics is the Bayes’ rule.

**Theorem 2.1.** Bayes’ rule. Assume we have defined probability distributions \( P(A) \), \( P(X|A) \), and \( P(X) \neq 0 \). Then

\[
P(A|X) = \frac{P(X|A)P(A)}{P(X)}. \tag{2.1}
\]

\( P(A) \) is the prior distribution for hypothesis \( A \). Typically, \( A \) is a hypothesis about some quantitative parameters \( \theta \). The prior measures the degree of belief about \( A \) being true before we have observed \( X \). \( P(X|A) \) is the likelihood of observing \( X \) given \( A \) is true. Evidence \( P(X) \) is the marginal probability of \( X \). The rule of total probability yields

\[
P(X) = P(X|A)P(A) + P(X|\neg A)P(\neg A).
\]
$P(A|X)$ is the posterior probability of $A$ given we have observed $X$. We can say that the Bayes’ rule updates the degree of the prior belief $P(A)$ yielding the updated belief as measured by the posterior probability $P(A|X)$.

The Bayes’ rule follows directly from the definition of conditional probability. A key characteristic of Bayesian inference in contrast with the frequentist paradigm is the use of probability for measuring uncertainty (Bernardo and Smith, 1994). In Bayesian inference, probabilities are interpreted as degrees of belief. In the frequentist interpretation, probabilities are considered as limits of relative frequencies. Especially, defining the prior distribution $P(A)$ with limits of relative frequencies can be problematic. Typically, it can be constructed with hypothetical repetitions under identical conditions (Gelman et al., 2004). But this can seem artificial, if for example we are modeling unique or almost unique events. Furthermore, in a strict frequentist interpretation an event that has not occurred would have probability zero, which can seem unintuitive in many situations (Hájek, 1996). In comparison, when $P(A)$ is interpreted as a prior belief about $A$, we get a coherent system of inference. The Bayesian approach to statistics can be derived from decision theoretic considerations (Bernardo and Smith, 1994). That way coherence of the inference system can be demonstrated. Alternatively, we can adopt the pragmatist view and state that Bayesian statistics has proven to be useful in the analysis of applied problems in many fields. This then justifies the use of (2.1) and the Bayesian approach in general.

### 2.2 Modeling

In practice, Bayesian inference includes but is not limited to specifying the likelihood and the prior. This is evident in Articles I, II and III which contain for example algorithms for learning the model in question. Defining these type of algorithms can be a central part of the statistical modeling effort.

When modeling real-world phenomena, all the relevant aspects of the process we are modeling cannot often be included in the model. There exists a separation between a ”theoretical world" and a "real world" (Kass, 2011). In the "theoretical world" we have mathematics and statistical models while
in the "real world" exist the actual phenomena to be studied and data. Statistical modeling can be seen as building a bridge between these two worlds. As part of the modeling, simplifying assumptions have to be made. Also, the usefulness of the model is something that has to be considered. A statistical model can be an adequate description of a phenomenon but not very useful if it is computationally too burdensome to be used for any purpose. This is evident for example when high-order Markov chains are used for analysing sequence data. Increasing the model order easily leads to models that are not useful because of the huge volume of computations involved.

One feature of the Bayesian approach is predictive inference. We can calculate marginal probability distributions for hypothetical and observed data. Also, predictive power is a useful measure of the suitability of a model (Bernardo and Smith, 1994). If the model predicts accurately new observations, then it seems to be an efficient description of the underlying phenomenon. A principle called Occam’s razor states that a simpler hypothesis should be preferred instead of unnecessarily complicated ones. In Bayesian statistics, this principle is at work in model comparison (MacKay, 1992). The idea is to calculate the marginal likelihood of the data for each competing hypothesis with (2.2). The probability distribution defined in (2.2) is also called the prior predictive distribution. These probabilities can be used for calculating Bayes factors (Kass and Raftery, 1995).

**Definition 2.2.** Marginal likelihood of the data and Bayes factor. For hypotheses, i.e. models $H_1$ and $H_2$ and data $X$, the marginal likelihood of data under the model $k, k = 1, 2$ is

$$P(X|H_k) = \int P(X|\theta_k, H_k)\pi(\theta_k|H_k)d\theta_k$$  \hspace{1cm} (2.2)

where $\theta_k$ denotes collectively parameters of the model $k$, $P(X|\theta_k, H_k)$ is the likelihood and $\pi(\theta_k|H_k)$ is the prior. Then the Bayes factor is

$$B_{12} = \frac{P(X|H_1)}{P(X|H_2)}.$$  \hspace{1cm} (2.3)
Here, the Bayes factor measures plausibility of $H_1$ in comparison to $H_2$. If the model has unnecessary parameters, this is penalized in the evidence (2.2) (MacKay, 1992).

A form of prediction is to calculate the posterior predictive distribution for future observations $\tilde{X}$ given data $X$ and a model. Following previous notation, the posterior predictive distribution of $\tilde{X}$ under model $k$ and given $X$ is

\[
P(\tilde{X}|X, H_k) = \int P(\tilde{X}|\theta_k, H_k)\pi(\theta_k|X, H_k)d\theta_k
\]

(2.4)

where $P(\tilde{X}|\theta_k, H_k)$ is the likelihood of $\tilde{X}$ and $\pi(\theta_k|X, H_k)$ is the posterior distribution of $\theta_k$.

In machine learning, data is often divided to subsets for learning and testing the model (Bishop, 2006). This is done to avoid over-fitting. If the model has a large number of parameters, it might describe well the training data but still fail to generalize to new data. The marginal likelihood of the data under a given model is a useful tool of Bayesian statistics that can be applied to many problems in machine learning. This is done in Articles I, II and III.
Markovian Models

Many fields of science and technology depend on the analysis of sequence data. Two prime examples are DNA sequences in biology and text in the context of processing natural language. Given empirical sequence data, a mathematical framework is needed for quantitative modeling (Koski, 2001, see also Article I). First, we consider a finite alphabet \( S = \{s_1, s_2, \ldots, s_J\} \) with \( J \) symbols. An example of a finite alphabet is the DNA alphabet \( \mathcal{X} = \{A, C, G, T\} \). We can label the symbols with integers and generally consider the following alphabet: \( \mathcal{X} = \{1, \ldots, J\} \). Let \( X_0, X_1, \ldots, X_n \) be a sequence of random variables that take values in \( \mathcal{X} \).

### 3.1 Markov Chain

A fundamental model for sequence data is the Markov chain.

**Definition 3.1.** Time homogeneous Markov chain (MC). Let \( \{X_n\}_{n=0}^{\infty} \) be a sequence of random variables. If for all \( n \geq 1 \) and \( j_0, j_1, \ldots, j_n \in \mathcal{X} \),

\[
P(X_n = j_n|X_{n-1} = j_{n-1}, \ldots, X_0 = j_0) = P(X_n = j_n|X_{n-1} = j_{n-1}),
\]

(3.1)

then \( \{X_n\}_{n=0}^{\infty} \) is called a Markov chain.

Elements in \( \mathcal{X} \) are called states of the Markov chain. A Markov chain has the *Markov property* defined by the equation (3.1). In essence, Markov property states that given the previous state \( X_{n-1} = j_{n-1} \), the rest of the history
is irrelevant for predicting the current state $X_n$. Probabilities of transitions between states can be represented in a transition matrix with elements $p_{ij} = P(X_n = j|X_{n-1} = i), i, j \in \mathcal{X}$. These probabilities are assumed to be independent of $n$ making the Markov chain *time homogeneous*.

Another important model is a Markov chain of $m$th order.

**Definition 3.2.** Time homogeneous Markov chain (MC) of $m$th order. Let $\{X_n\}_{n=0}^{\infty}$ be a sequence of random variables. If for all $n \geq m$ and $j_0, j_1, \ldots, j_n \in \mathcal{X}$,

$$P(X_n = j_n|X_{n-1} = j_{n-1}, \ldots, X_0 = j_0) = P(X_n = j_n|X_{n-1} = j_{n-1}, \ldots, X_{n-m} = j_{n-m}),$$

for a positive integer $m$, then $\{X_n\}_{n=0}^{\infty}$ is called a Markov chain of $m$th order.

A Markov chain is called a first-order Markov chain in this context. A Markov chain $\{X_n\}_{n=0}^{\infty}$ of $m$th order can be transformed to a first-order Markov chain $\{Z_n\}_{n=0}^{\infty}$ with transition probabilities $p_{ij}, i, j \in \mathcal{X}^m$ and an extended state space with $|\mathcal{X}|^m = J^*$ states. Transition matrix of the first-order Markov chain then has $|\mathcal{X}|^m$ rows and each row has exactly $J$ transition probabilities that can have positive values.

Markov chains are useful for modeling different type of phenomena. Often, the Markov assumption describes reasonably well the dynamics of a real-world system. Heuristically, it is plausible to assume that what is close proximity affects the current state more than what is further apart.

In a Markov chain model, the order of the model controls how much of the history of the process is used at each step. In principle, it would be tempting to examine all of the history at once. However, realities of modeling limit the possibilities. The number of free parameters for a Markov chain of $m$th order is $|\mathcal{X}|^m(|\mathcal{X}| - 1)$. A Markovian model that has this number of free parameters is called a full Markov chain. This number grows exponentially with the order of the Markov chain. Thus, the number of observed transitions per state decreases when the order grows. This leads to difficulties in the estimation of transition probabilities. Also, this means
that finding an optimal order $m$ for the Markov chain model is far from trivial.

Due to challenges associated with using $m$th order Markov chains, several alternatives have been presented. They are typically based on the idea of sparsity, aiming to find parsimonious descriptions of the data generating process. Starting from the $m$th order Markov chain model, we can try to find models that utilize the same information but in a more effective manner. This leads to the question of data compression. There is a rich literature on the Variable order Markov models (VOM) starting from Rissanen (1983). Here we focus on Variable length Markov chains, Sparse Markov chains and Mixture transition distribution models.

### 3.2 Variable Length Markov Chain

A realized value for a random variable $X_t$ is denoted by $x_t$.

**Definition 3.3.** Variable length Markov chain (VLMC). Let $\{X_n\}_{n=0}^\infty$ be a time homogeneous Markov chain of $m$th order. Denote by $c_{\text{pre}} : \mathcal{X}^m \to \bigcup_{j=0}^m \mathcal{X}^j$ a function which maps $c_{\text{pre}} : x_{n-1}, \ldots, x_{n-m} \mapsto x_{n-1}, \ldots, x_{n-l}$ where

$$l = l(x_{n-1}, \ldots, x_{n-m}) = \min\{k \in \mathbb{Z}_{\geq 0} \mid P(X_n = j_n | X_{n-1} = x_{n-1}, \ldots, X_{n-m} = x_{n-m}) = P(X_n = j_n | X_{n-1} = x_{n-1}, \ldots, X_{n-k} = x_{n-k}) \text{ for all } j_n \in \mathcal{X}\},$$

such that $l = 0$ corresponds to independence. Then $l$ is a variable length memory and $c_{\text{pre}}(\cdot)$ is the preliminary context function. Final context function $c(\cdot)$ is obtained by lumping together some of the values of $c_{\text{pre}}(\cdot)$ that share the second to last symbol. A Markov chain of $m$th order with a variable length memory $l$ is called a Variable length Markov chain of order $p$ where $p$ is the smallest integer such that $l(x_{n-1}, \ldots, x_{n-m}) \leq p \leq m$ for all $x_{n-1}, \ldots, x_{n-m} \in \mathcal{X}^m$.

This definition of Variable length Markov chain is slightly different from those in Bühlmann and Wyner (1999) and Mächler and Bühlmann (2004).
as they start from a stationary process in finite alphabet instead of a Markov chain of \( m \)th order as is done here.

### 3.3 Mixture Transition Distribution

Another parsimonious model for high-order Markov chains is the Mixture transition distribution model (Raftery, 1985; Le et al., 1996; Berchtold and Raftery, 2002). It has been extended from modeling of high-order Markov chains in a finite state space to for example general state spaces. Here, the focus is on applying MTD models for high-order Markov chains in a finite state space.

**Definition 3.4.** Mixture transition distribution (MTD). Let \( \{X_n\}_{n=0}^{\infty} \) be a time homogeneous Markov chain of \( m \)th order. In the corresponding MTD model we have

\[
P(X_n = j_n | X_{n-1} = j_{n-1}, \ldots, X_{n-m} = j_{n-m}) = \sum_{g=1}^{m} \lambda_g P(X_n = j_n | X_{n-g} = j_{n-g}) = \sum_{g=1}^{m} \lambda_g p_{j_n|j_{n-g}}
\]

with constraints

\[
\sum_{g=1}^{m} \lambda_g = 1, \quad \lambda_g \geq 0.
\]

In the MTD model, contributions of the lags \((1 \ldots m)\) are combined additively. This offers a parsimonious model of the Markov chain. The model has \(|\mathcal{X}|(|\mathcal{X}|-1)+(m-1)\) free parameters which is clearly less than \(|\mathcal{X}|^m(|\mathcal{X}|-1)\) of the full \( m \)th order Markov chain.
3.4. SPARSE MARKOV CHAIN

VLMC and MTD are complementary models in a sense that they have
different strengths (Berchtold and Raftery, 2002). MTD model of $m$th order
always deals with $m$ lags while in the VLMC model part of the history is
irrelevant depending on the context. Berchtold and Raftery report a com-
parison of these two models.

3.4 Sparse Markov Chain

Finally, we present the Sparse Markov chain (SMC). Such models are defined
in Article I, and they offer another alternative to $m$th order Markov chains.

Definition 3.5. Sparse Markov chain (SMC). Let $\{X_n\}_{n=0}^{\infty}$ be a time ho-
moregous Markov chain of finite order $m$ transformed to a first-order MC
$\{Z_n\}_{n=0}^{\infty}$. Let $S = (s_1, \ldots, s_k)$ be a partition of $X^m$ such that the transi-
tion probability vectors satisfy the equality $p_{ij} = p_{j}$. for all pairs of states
$\{i, j\} \in s_c, c = 1, \ldots, k$, and $P$ the corresponding set of $k$ transition prob-
ability distributions in $X^m$. If $k < |X^m|$, the pair $(S, P)$ is called an SMC
(of order $m$).

The following theorem characterizes connection between SMC and VLMC
models.

Theorem 3.6. Let $(S, P)$ be an SMC. Then, there is an equivalent repre-
sentation based on the set of contexts $\mathcal{B}$ of a VLMC model if and only if
there exists a unique context $B_c$ with $b^{(r)}$, which is a suffix to all states $i$
assigned to the same class $s_c$ for all $c = 1, \ldots, k$.

The proof is given in Article I. Theorem (3.6) formally states that for
a VLMC model that is not a full Markov chain, there is an equivalent
representation as an SMC model. The reverse is not generally true. There
are SMC models that do not have representation as a VLMC model as is
shown in the proof. For example, in Article I an SMC model is described
for which $X_{n-1}$ is irrelevant for predicting $X_n$ while $X_{n-2}$ is relevant. This
kind of probability model cannot be described with the context tree of a
VLMC model.
For estimating Markovian models, there are plenty of methods. For any Markov chain of order $m$, maximum likelihood based methods can be used. For VLMC models, context algorithm can be used (Bühlmann and Wyner, 1999; Rissanen, 1983). For a MTD model, numerical maximization of the log-likelihood or the expectation-maximization (EM) algorithm can be used (Berchtold and Raftery, 2002). In Chapter 5, a Bayesian method for learning Markovian models is presented.
Graphical Representations

Earlier in Chapter 3, the following Markovian models were introduced: MC, VLMC and SMC. Here, we give a graphical representation of these models (see Article I). We find a DAG (directed acyclical graph) for the sample paths of each of the models. In all of them, we have a tree structure that represents probability distributions for the random variable $X_n$ conditional on values of $X_{n-1}, \ldots, X_{n-m}$. Generally, a DAG consists of nodes and directed edges between them so that there are no cycles (Koski and Noble, 2009). Here, $X_n$ corresponds to the root node in the graph. Other nodes in the tree consist of possible values (or a set of them) of $X_{n-1}, \ldots, X_{n-m}$. Here, $m$ is the order of the Markovian model. Directions of edges are from $X_{n-m}$ to $X_n$.

Firstly, we consider a full MC of order $m$ in DNA alphabet. All the possible histories are represented in a tree. Example is given in Figure 4.1. The number of leaves is $|\mathcal{X}|^m$ while the number of free parameters is $|\mathcal{X}|^m(|\mathcal{X}| - 1)$. Here, leaves are those nodes for which no edges are directed at.

For a VLMC model of order $m$, part of the full tree has been pruned. This corresponds to the situation where for given a set of sample paths the same probability distribution always holds for the random variable $X_n$. In comparison with the SMC model, this set of sample paths has to be hierarchical, i.e. it has to have one common path to the root node $X_n$. Nodes can be lumped together in two ways. Firstly, history beyond some node can be irrelevant. For example, in Figure 4.2 if $X_{n-1} = T$, then all
values of $X_{n-2}$ lead to the same probability distribution for $X_n$. Secondly, nodes sharing a second to last symbol can result in the same probability distribution for $X_n$. In Figure 4.2, histories $CC$, $CG$ and $CT$ are lumped together.

Also, for a VLMC model there exists an SMC representation which is a partition of the sample paths of the full tree of order $m$. For the VLMC model in Figure 4.2, there would be 11 clusters corresponding with the leaves of the pruned tree.

It should be also noted that the VLMC model’s dependence structure
in the space of sample paths corresponds to the theory of context-specific DAGs, namely labeled DAGs, as presented by Pensar et al. (2014).

For an SMC model, it is possible that there does not exist a corresponding VLMC model. This is illustrated in Figures 4.3 and 4.4. Numbers inside the nodes denote to which cluster the sample paths belong to.

In Figure 4.3, we have the previously mentioned example, where $X_{n-1}$ is irrelevant for predicting $X_n$ while $X_{n-2}$ is relevant. There are four clusters in the partition of the sample paths of the full three. In Figure 4.4, we have three clusters and no clear hierarchical structure as generally nodes in one cluster do not share a path to the root node $X_n$. In learning a Markovian model, we can estimate the partition of a SMC model. This enables the learning of MC and VLMC as well as SMC models.
Figure 4.4: A sparse Markov chain of order 2 in DNA alphabet with 3 clusters
Bayesian Clustering

Here, we describe a Bayesian approach to clustering. In Articles I, II and III, Bayesian clustering is used for learning sequence models from data. However, there are certain differences between the considered problems. For example, in Articles I and III states of Markovian models are clustered while in Article II the objects to be clustered are sequences modeled with Markov chains of fixed order. We aim to present the general ideas of Bayesian clustering as well as to elaborate on some important details.

In Article II, we define a partition $S$ of a set $S_0$ as "a collection of disjoint, non-empty subsets of $S_0$, whose union is $S_0". The elements of the partition are called clusters. Clustering means finding a partition following some criterion. Typically, this is similarity. Then there should be a high probability that similar objects belong to the same cluster. Clustering is an example of a task in unsupervised machine learning (Bishop, 2006). The procedure is unsupervised because in principle no information besides the set of data objects is used.

### 5.1 Prior and Likelihood

In Bayesian clustering, the partition $S$ is the variable of main interest. A central idea is that the partition $S$ with high posterior probability $p(S|x)$ should be close to an optimal partition. This can be justified with results on predictive inference and classification theory (Corander et al., 2007, 2013;
Here, $x$ denotes a set of observations with $n$ objects to be clustered. To calculate the posterior distribution $p(S|x)$, a prior $p(S)$ and a marginal likelihood $p(x|S)$ are needed. Then the Bayes’ rule (2.1) yields

$$p(S|x) = \frac{p(x|S)p(S)}{\sum_{S \in S} p(x|S)p(S)},$$

(5.1)

where $S$ is the set of all possible partitions of $x$. In the expression for the marginal likelihood $p(x|S)$, parameters of the model have been integrated over their prior distributions. A reasonable approach to identifying a good partition is to solve the maximum a posteriori (MAP) estimate of the partition parameter $S$. The MAP estimate is defined as

$$\hat{S} = \arg\max_{S \in \mathcal{S}} p(S|x)$$

$$= \arg\max_{S \in \mathcal{S}} p(x|S)p(S).$$

(5.2)

Often, the MAP estimate can be solved only approximately. A simple solution to finding the MAP estimate would be to evaluate $p(x|S)p(S)$ for all possible values of $S$. However, this is computationally infeasible with almost any realistic data set. The number of possible partitions for $n$ objects is the Bell number $B(n)$ and it increases rapidly as a function of $n$ (Bell, 1934; Rota, 1964). There are various methods that a stochastic algorithm could employ to maximize $p(x|S)p(S)$. Markov chain Monte Carlo (MCMC) and similar stochastic simulation methods give consistent MAP estimates but they can be too slow in some cases when the number of data objects to be clustered is large. In Articles I and II, a stochastic greedy algorithm is used. Search operators like joining two clusters together and splitting one cluster into two are used in a data-driven manner. In Article III, a deterministic recursive learning algorithm is used. In later sections, algorithms for learning the MAP partition are presented.

The prior $p(S)$ can take many forms. The simplest form is the uniform prior $p(S) = 1/B(n)$. Then each partition is equally probable and it is
enough to maximize $p(x|S)$ when calculating the MAP estimate. Another possibility is the Dirichlet process prior (DPP)

$$p(S) \propto \prod_{S_c} \eta_0 \Gamma(|S_c|), \quad (5.3)$$

where $\eta_0$ is a hyperparameter and $|S_c|$ is the number of items in cluster $S_c$. Here, $\Gamma(\cdot)$ denotes the Gamma function. With the Dirichlet process prior, sizes of individual clusters affect the prior probability of the partition so that partitions with larger clusters are favoured. This is desirable if we have a priori information that the number of clusters should be small relative to $n$, the number of items to be clustered.

The following description of the Dirichlet process prior is adapted from Article II. To derive the Dirichlet process prior (5.3) we can consider the Dirichlet process mixture (DPM) model under certain assumptions (Neal, 2000; Jain and Neal, 2004; Teh et al., 2006; Dahl, 2009). By representing the Dirichlet process (Ferguson, 1973) with the Pólya urn scheme (Blackwell and MacQueen, 1973), the prior can derived (see, e.g. Dahl, 2009). The stick-breaking construction of the Dirichlet process (Sethuraman, 1994) provides an alternative approach. Then we have the following form of the DPM model: (see, e.g. Teh et al., 2006)

$$\pi | \eta_0 \sim \text{GEM}(\eta_0)$$
$$z_i | \pi \sim \pi \quad (5.4)$$
$$\varphi_k | G_0 \sim G_0$$
$$y_i | z_i, (\varphi_k)_{k=1}^{\infty} \sim F(\varphi_{z_i}),$$

where given the cluster membership indicators $z_i$ and cluster parameters $\theta_i = \varphi_{z_i}$ we have conditional independence of observations $y_i$. The prior distribution for cluster parameters is $G_0$. We draw indicator variables $z_i$ independently from the stick-breaking distribution $\pi$ which is sometimes denoted by GEM($\eta_0$). GEM stands for Griffiths, Engen, and McCloskey.
The stick-breaking construction of the Dirichlet process DP($\eta_0, G_0$) induces the random probability measure on the positive integers $\pi = (\pi_k)_{k=1}^{\infty}$. We have $G = \sum_{k=1}^{\infty} \pi_k \delta_{\varphi_k}$ in the construction. The distribution of $G$ follows DP($\eta_0, G_0$) (Sethuraman, 1994). By assuming $n$ observations from the DPM model above, we get a finite number of clusters that have more than zero observations. We have a multinomial distribution with event probabilities defined by $\pi$ and the cluster sizes follow this distribution. In the stick-breaking construction, there is a beta distributed variable $\beta_k$ associated with each $\pi_k$ (see, e.g. Teh et al., 2006). We can use the Bayes’ theorem to get a posterior distribution for the cluster sizes and the unknown $(\beta_k)_{k=1}^{\infty}$. By integrating out each $\beta_k$ over its beta prior distribution we get the Dirichlet process prior as given in (5.3).

Typically, in Bayesian clustering it is assumed that data objects in different clusters are conditionally independent given the partition. Then, we have a product partition model (Hartigan, 1990; Barry and Hartigan, 1992). In these models, the likelihood is expressed as a product

$$p(x|S) = \prod_{S_c \in S} f(x_{S_c}), \quad (5.5)$$

where $f(x_{S_c})$ is the marginal likelihood contribution from cluster $S_c$, which can take a variety of forms. Here, the likelihood contribution from a cluster $S_c$ is defined as

$$f(x_{S_c}) = \int_{\Theta} p(\theta)p(x_{S_c}|\theta)d\theta, \quad (5.6)$$

where $\theta$ denotes collectively the quantitative parameters of the model.

Next, we define the marginal likelihood $p(x|S)$ for Markovian models excluding MTD. The marginal likelihood for SMC models is a basis for the learning algorithms. Because a VLMC model has a representation as an SMC model, the same definition for the marginal likelihood can be used when learning VLMC and SMC models. Also, the marginal likelihood for a full MC can be calculated with the same formulations.
5.1. PRIOR AND LIKELIHOOD

For an MC of given order $m$, we have data on transitions from $|\mathcal{X}^m|$ states to $J = |\mathcal{X}|$ symbols. For an SMC model with a partition $S$ and $k$ clusters, there are $\{p_c\} : c = 1, \ldots, k$ probability vectors. For a full MC, $k = |\mathcal{X}^m|$. We denote by $\theta \in \Theta$ the quantitative parameters of the model. When assuming the initial state fixed, the likelihood is a product of multinomial distributions

$$p(x|\theta, S) \propto \prod_{i=1}^{|\mathcal{X}^m|} \prod_{j=1}^J n_{ij}^{p_{c_{ij}}} = \prod_{c=1}^k \prod_{j=1}^J p_{c_{ij}}^{\sum_{i \in sc} n_{ij}}, \quad (5.7)$$

where the number of transitions from the state $i$ to $j$ in $x$ is denoted by $n_{ij}$. For transition probabilities $p_{c_{ij}}, c = 1, \ldots, k, j = 1, \ldots, J$, we choose the canonical conjugate multivariate Dirichlet prior (Koski, 2001)

$$p(\theta|\alpha, q) = \prod_{c=1}^k \left[ \frac{\Gamma(\alpha)}{\prod_{j=1}^J \Gamma(\alpha q_j)} \prod_{j=1}^J p_{c_{ij}}^{\alpha q_j-1} \right], \quad (5.8)$$

where the hyperparameters satisfy the following conditions: $\alpha > 0, q_j > 0, \sum_{j=1}^J q_j = 1$. Using the properties of Dirichlet distribution, the marginal likelihood $p(x|S)$ can be calculated as

$$p(x|S) \propto \int_{\theta \in \Theta} p(x|\theta, S)p(\theta|\alpha, q) d\theta$$

$$\propto \int_{\theta \in \Theta} \left[ \prod_{c=1}^k \prod_{j=1}^J \frac{\Gamma(\alpha)}{\Gamma(\alpha q_j)} p_{c_{ij}}^{\alpha q_j-1} \prod_{j=1}^J p_{c_{ij}}^{\sum_{i \in sc} n_{ij}} \right] d\theta \quad (5.9)$$

$$\propto \prod_{c=1}^k \frac{\Gamma(\alpha)}{\prod_{j=1}^J \Gamma(\alpha q_j)} \prod_{j=1}^J \frac{\Gamma(\sum_{i \in sc} n_{ij} + \alpha q_j)}{\Gamma((\sum_{j=1}^J \sum_{i \in sc} n_{ij}) + \alpha)}.$$
One important issue is choosing $m$, the order of the Markov chain used in estimation. In Article I, a uniform distribution is assigned over the values $m = 0, \ldots, M$, where $M$ is an upper bound that has to be chosen and can be revised if necessary. Then the joint posterior distribution of $m$ and $S$ is

$$p(S, m \mid x) \propto p(x \mid S)p(S)p(m)$$

and $p(m) = 1/(M + 1)$. The approximate MAP estimate can be calculated for both $S$ and $m$:

$$\hat{S}, \hat{m} = \arg\max_{m \in \{0, \ldots, M\}} \left\{ \arg\max_{S_m \in S_m} p(x \mid S_m)p(S_m) \right\}.$$ (5.11)

### 5.2 Expectation-Maximization Algorithm

Here, we discuss the expectation-maximization (EM) algorithm. Its first general formulation is due to Dempster et al. (1977). In Article II the EM-algorithm is used for learning a partition when the data objects are Markov chains of fixed order. In the EM-algorithm, we have unknown and latent variables. For example, a latent variable can be the partition $S$ while the unknown variables can be the transition probability matrices collectively denoted by $\theta$. Assume that $x$ represents the observed data, $z$ is a latent variable and $\theta$ is an unknown parameter. Then consider a posterior $p(\theta \mid x) = \sum_z p(\theta, z \mid x)$. Here, the description of the algorithm is adapted from Article II. The EM-algorithm for finding a MAP estimate of the posterior is then defined as follows. Starting with an initial value $\theta^{(0)}$ for $\theta$ and setting $k \leftarrow 0$, the following steps are applied until convergence:

- **Expectation step**
  
  Calculate

  $$Q(\theta \mid \theta^{(k)}) = E_{\theta^{(k)}, x} \left[ \ln p(\theta, Z \mid x) \right]$$

  $$= \sum_z \ln p(\theta, z \mid x) \cdot p(z \mid \theta^{(k)}, x)$$


5.3. **CLASSIFICATION EM ALGORITHM**

- **Maximization step**
  
  Set 
  
  \[
  \theta^{(k+1)} \leftarrow \arg \max_{\theta} Q(\theta|\theta^{(k)})
  \]
  
  and \(k \leftarrow k + 1\).

  Convergence is achieved when the difference between \(\theta^{(k+1)}\) and \(\theta^{(k)}\) is below a threshold that has been set beforehand. Also, it can be useful to define an upper bound for the number of iterations that the algorithm is allowed to run.

  In Article II, the EM-algorithm is used for estimating transition probability matrices of sequences modeled with Markov chains and for assigning sequences to the most appropriate clusters. The number of clusters is not changed by the EM-algorithm.

  For iterations \(k\) and \(k+1\) it holds that 
  
  \[
  p(\theta^{(k+1)}|x) \geq p(\theta^{(k)}|x).
  \]

  Thus, the definition of the EM-algorithm leads to monotonic probabilities for the unknown parameter \(\theta\) given the data \(x\). This property of the algorithm makes the convergence possible. Here, we have marginalized out the latent variable \(z\).

### 5.3 Classification EM Algorithm

In Article IV, we have \(M\) generating clusters (mechanisms) and \(N\) observations. The aim is then to estimate the prior probability that an observation is generated by cluster \(j\), denoted by \(\alpha_j\) for \(j = 1, \ldots, M\). Also, we estimate parameter vectors for mechanisms \(j = 1, \ldots, M\), denoted by \(\psi_j\). Finally, assignment of each observation to a cluster, denoted by \(u^{(l)}\) for \(l = 1, \ldots, N\) and \(j = 1, \ldots, M\), has to be estimated. For an alternative representation, indicator variables of the form \(u_{j}^{(l)}\) can be used. In article IV, a version of classification EM algorithm (CEM) (Celeux and Govaert, 1992; Redner and Walker, 1984) is used. The following description of the classification EM algorithm is adapted from Article IV:
• Initial step
  Choose initial values for $\alpha_j^{(0)}$ and $\psi_j^{(0)}$, $j = 1, \ldots, M$. Move to E-step.

• First step
  $\alpha_j^{(0)}$ and $\psi_j^{(0)}$, $j = 1, \ldots, M$ are our current parameters and $u_j^{(l)}$, $l = 1, \ldots, N$ are our current assignments.

• E-step
  Compute for $l = 1, \ldots, N$ and $j = 1, \ldots, M$ using the Bayes’ rule

  $$t_j(y^{(l)}) = \frac{\alpha_j^{(0)} p(y^{(l)}|\psi_j^{(0)}, j, x^{(l)})}{\sum_{j=1}^{M} \alpha_j^{(0)} p(y^{(l)}|\psi_j^{(0)}, j, x^{(l)})}.$$  

  Here, $t_j(y^{(l)})$ is the current posterior probability that $(y^{(l)}, x^{(l)})$ belongs to the cluster $j$.

• C-step
  For $l = 1, \ldots, N$ we assign

  $$\hat{u}_j^{(l)} = \begin{cases} 1 & j^* = \arg\max_{1 \leq j \leq M} t_j(y^{(l)}) \\ 0 & \text{otherwise} \end{cases}$$

  to get a new assignment of $(y^{(l)}, x^{(l)})$. Thus, each observation is assigned to a cluster so that the posterior probability is maximized.

• M-step

  $$\alpha_j^{(1)} = \frac{\sum_{l=1}^{N} u_j^{(l)}}{N}$$
5.4. **STOCHASTIC SEARCH**

is maximum likelihood estimate of $\alpha_j$ from $L_2(\alpha) = \sum_{j=1}^{M} n_j \log(\alpha_j)$.

We need to find $\psi_j^{(1)}$ for $j = 1, \ldots, M$ by maximization of $L_1(\psi) = \sum_{j=1}^{M} L_j(\psi_j)$. Details of this maximum likelihood estimation are given in Article IV.

- **Return to**

Now, we have the new estimates $\alpha_j^{(1)}, \psi_j^{(1)}, u_l^{(l)}$, $j = 1, \ldots, M; l = 1, \ldots, N$.

We assign $\alpha_j^{(1)} \rightarrow \alpha_j^{(0)}, \psi_j^{(1)} \rightarrow \psi_j^{(0)}, u_l^{(l)} \rightarrow u_l^{(l)}$.

- **Stop**

We have $\ln(L) = L_1(\psi) + L_2(\alpha)$. The algorithm is stopped when

$$|\ln L(\psi^{(1)}, \alpha^{(1)}) - \ln L(\psi^{(0)}, \alpha^{(0)})| \leq \varepsilon$$

where $\varepsilon$ is a small positive number that has to be specified in implementation. Also, if a preassigned maximum number of iterations has been reached, the algorithm stops.

### 5.4 Stochastic Search

Here, we present a stochastic search algorithm adapted from Article I. This greedy algorithm is based on algorithms presented in Corander and Marttinen (2006) and Marttinen et al. (2006). The idea is to find a good clustering of the $|X|^m$ states of the Markov chain and thus estimate the Sparse Markov chain model. For a given value of $m$ we have the following algorithm:

(i) Initialize $S_t, t = 0$ with $|X|^m$ singleton clusters and store for all pairs of states $i, l \in X^m$ the distances between posterior mean estimates of their transition probability vectors.
\[ d_{i,l} = \sum_{j=1}^{J} \left( \frac{n_{i|j} + \alpha q_j}{\sum_{j=1}^{J} n_{i|j} + \alpha q_j} - \frac{n_{l|j} + \alpha q_j}{\sum_{j=1}^{J} n_{l|j} + \alpha q_j} \right)^2. \] (5.12)

(ii) Given the current value of \( p(x|S_t) \), apply the following operators sequentially until no change in \( S_t \) results in a higher marginal likelihood.

(iii) In a random order, move each state \( i \in X_m \) to the class \( c \in S_t \), which results in the \( S_{t+1} \) associated with a maximal increase in \( p(x|S_{t+1}) \). If \( p(x|S_{t+1}) \leq p(x|S_t) \) for all \( c = 1, \ldots, k, S_{t+1} = S_t \).

(iv) For each pair of classes \( c, c' = 1, \ldots, k \), calculate \( p(x|S^*) \) for the \( S^* \) obtained by merging classes \( c, c' \) in \( S_t \). If any \( S^* \) satisfies \( p(x|S^*) - p(x|S_t) > 0 \), set \( S_{t+1} \) equal to the \( S^* \) for which \( p(x|S^*) - p(x|S_t) \) is maximal, otherwise set \( S_{t+1} = S_t \).

(v) For each class \( c = 1, \ldots, k \), use the complete linkage algorithm (e.g. Mardia et al., 1979) with distances (5.12) to split the class into two non-empty subsets of states and calculate \( p(x|S^*) \) for the resulting partition \( S^* \). If \( p(x|S^*) - p(x|S_t) > 0 \), set \( S_{t+1} \) equal to \( S^* \), otherwise set \( S_{t+1} = S_t \).

This greedy algorithm converges to a local mode when the operators do not increase marginal likelihood any further. Several restarts from different initial conditions can be used to find clusterings that are closer to a global optimum. Also, this algorithm could be generalized to give a consistent posterior estimator using Markov chain Monte Carlo approach with a non-reversible Markov chain (Marttinen et al., 2006; Corander et al., 2006, 2008).

5.5 Recursive Search

In Article III, a heuristic deterministic algorithm is defined for searching the optimal SMC model \((S, P)\) for a given sequence \(\{X_t\}_{t=1}^n\). Here, we present an
adapted description of the algorithm. The general idea is to apply Delaunay triangulation on $\mathcal{X}^m$ so that each transition state $i \in \mathcal{X}^m$ becomes a node in the triangulation graph. Then we recursively merge nodes to maximize the posterior probability. Bayes factor is used as the local criterion when choosing which nodes should be merged.

- **Initial step**
  Obtain the transition counts of $\{X_t\}_{t=1}^n$ for MC of order $m$. Estimate the transition probability distribution $\theta$ of the model by posterior mean estimation from the transition counts. Form Delaunay triangulation $G$ of $\mathcal{X}^m$ by using values of free parameters in $\theta$ as coordinates. Calculate the log Bayes factor $\log BF_{uv}$ for each edge $u, v$ in $G$. Find the edge $(u^*, v^*)$ that gives maximal log Bayes factor value $w$. Set $\mathcal{U} = u^*, \mathcal{V} = v^*$ and $\mathcal{W} = w$.

- **While $\mathcal{W} > 0$ do**
  Merge $\mathcal{V}$ to $\mathcal{U}$ by the following steps:
  a) add the sufficient statistics counts of $\mathcal{V}$ to $\mathcal{U}$
  b) for each node $r$ in $G$ which has a connection with $\mathcal{V}$, if edge $(\mathcal{U}, r)$ does not exist, redirect the edge $(\mathcal{V}, r)$ to $(\mathcal{U}, r)$
  c) delete $\mathcal{V}$ from $G$. Update the Bayes factors for all the edges (include the edges added by merging) connected to $\mathcal{U}$. Find the edge $(u^{**}, v^{**})$ with a maximal log Bayes factor value $w'$. Set $\mathcal{U} = u^{**}, \mathcal{V} = v^{**}$ and $\mathcal{W} = w'$. 
Causal inference is a challenging issue for both scientific inquiry and philosophy of science (Rothman et al., 2008). In this chapter, we present two models of causality. The first one is Rothman’s pie model, also called the sufficient-component cause model (Rothman, 1976; Rothman et al., 2008; Rothman, 2012). The second one is Pearl’s Do-Calculus which gives a probabilistic account of causality (Pearl, 1995, 2000). In Article IV, a model that is a stochastic extension of the sufficient-component cause model is presented. In the article, we have a fixed number of causal mechanisms. The effect of the covariates on the probability distribution of the outcome variable depends on the mechanism-specific parameters. These are estimated from the data so that mechanisms can be identified. We can consider mechanisms analogous to sufficient causes in the sufficient-component cause model. When simulating data from the model, we first draw randomly a sufficient cause and then based on that we draw the value of the outcome variable. In this model, the individual component causes are modeled with parameters of the mechanisms.

There is a connection between the model in Article IV and Pearl’s Do-Calculus. Also, Do-Calculus is described here as an example of a probabilistic framework for causality. Because the model in Article IV is essentially a Bayesian network, we can apply do-conditioning. The covariates do not have parents in the graph so forcing a subset of covariates to have certain value leads to a equivalent conditional probability distribution as see-conditioning, i.e. observing those values.
Traditionally, statistical inference has been concerned with associations and correlations instead of causality. In scientific inquiry, causal relations are important and often results of statistical methods are used for causal inference, with mixed success (Freedman, 1999). However, in statistics there are several important and relatively successful approaches to causal inference. Firstly, there is counterfactual causality (see, e.g. Rubin, 1974). For example, in a typical medical study each subject either gets the treatment or does not. We are then interested in the causal effect, i.e. the difference in outcome between a situation where the subject gets the treatment and the alternative situation where the subject does not get the treatment. Here, counterfactuality comes from the fact that the subject cannot both get the treatment and not get it. Randomization makes it possible to assess the causal effect when there are several trials. Rubin also discusses how an observational, non-randomized study can provide information on the causal effect. The counterfactual or potential outcomes approach has been extended for example with the use of structural equations and instrumental variables (Angrist et al., 1996). Another approach is the use of graphical models for causal inference (see, e.g. Pearl, 2000). Often, relations between variables are described with directed acyclic graphs (DAGs) and parents of a node are considered to be its direct causes. There has been symbiotic development between the counterfactual and graphical models approaches (Greenland and Brumback, 2002; Pearl, 2009). Finally, there is the predictive causality approach which takes explicitly into account the time between the cause and the effect (see, e.g. Arjas and Eerola, 1993; Arjas and Parner, 2004).

6.1 Sufficient-Component Cause Model

In epidemiology, causation is often modeled with the sufficient-component cause model. Here, we state the basic principles of the model (see e.g. Rothman et al., 2008). We define cause as a condition or event that precedes the disease and had the cause not been present the disease would not have occurred. By sufficient cause we mean a set of component causes that is minimal and complete and which is sufficient for the disease to appear.
6.1. SUFFICIENT-COMPONENT CAUSE MODEL

Minimality means that if one of the component causes is not present than the disease does not appear. Completeness means that if all the component causes are present than the disease occurs. There are typically several different sufficient causes for a disease. If a component cause is present in all sufficient causes, we can consider it to be necessary.

These principles are illustrated in Figure 6.1. There we have component causes $A$, $E$, $B$ and $D$. These are assumed to be binary so that the condition of the component cause either is present (value 1) or not (value 0). Typically, in epidemiology there are always component causes that have not been identified. This is taken into consideration by having an unknown cause $U_i$ as a part of each sufficient cause. We notice from the figure that the value $B = 0$ is part of the first sufficient cause while the value $B = 1$ is part of the third sufficient cause. This could be the situation for example if presence of some chemical substance would in one case be causing a disease and in another case preventing it. This depends on the other component causes of the sufficient cause. They are called collectively causal complement. Thus, the effect of chemical substance being present or not depends on its causal complement.

One insight that the sufficient-component cause framework gives is that the proportion of disease due to specific causes can add to over 100%. This is because there are multiple sufficient causes that have at least partially different component causes.

The sufficient-component cause model is deterministic but the methodology of epidemiology includes a variety of mathematical tools (Rothman et al. 2008).
et al., 2008; Rothman, 2012). For example, the risk of contracting a disease for an individual can be modeled with probabilistic methods. This modeling can benefit from the sufficient-component cause model when different causes are assessed systematically. Also, a stochastic generalization of the sufficient-component cause model can be useful for epidemiological study.

It seems that for causal inference, a logical, qualitative foundation is needed. The sufficient-component cause model provides this. Then causal inference can proceed with either a model that includes causality or with more heuristic methods, perhaps comparing information from many sources or experiments. Article IV is an example of work that develops a probabilistic model of causality.

6.2 Do-Calculus

A notable framework for probabilistic causal inference is Pearl’s Do-Calculus (Pearl, 1995, 2000). We present basic definitions and some properties of Do-Calculus following Koski and Noble (2009) as well as lecture notes by Koski and Noble and lecture slides from Koski’s presentation from the Bayesian network course at KTH, year 2013.

Firstly, we define a Bayesian network as the following structure.

**Definition 6.1.** Bayesian Network. Let $G$ be a directed acyclic graph (DAG) with nodes $V = \{1, \ldots, n\}$, $(X_v)_{v \in V}$ a set of finite discrete random variables and $\{P(X_v = x_i | X_{\pi(v)} = x_{\pi(v)})\}_{v \in V}$ a set of conditional probability distributions with

$$P(X_1 = x_{i_1}, \ldots, X_n = x_{i_n}) = \prod_{v=1}^{n} P(X_v = x_{i_v} | X_{\pi(v)} = x_{\pi(v)}),$$

where $\pi(v)$ is the set of parent nodes of $v$ and $P(X_v | X_{\pi(v)}) = P(X_v)$ if $\pi(v) = \emptyset$.

We consider the parents $X_{\pi(v)}$ as the direct causes of $X_v$ and $P(X_v | X_{\pi(v)})$ measures our belief about the strength of the causality. Often, in statistical
inference we observe a random variable having some value. In controlled experiments, a variable is forced to have some value (maybe after randomization). A central part of Do-Calculus is the intervention formula. It describes the joint probability distribution of the Bayesian network after a subset of variables has been forced some values. Observing a value is called see-conditioning and forcing a value is called do-conditioning.

When variables $X_A$, $A \subseteq V$ are forced to have values $x_A^*$ the resulting joint distribution is defined as

$$P(X_{V} = x_{V}|X_{A} \leftarrow x_{A}^*) = \frac{P(X_V = x_V)}{\prod_{v \in A} P(X_v = x_{i_v}|X_{\pi(v)} = x_{\pi(v)}) | x_A = x_A^*}$$

$$= \prod_{v \in V \setminus A} P(X_v = x_{i_v}|X_{\pi(v)} = x_{\pi(v)}) | x_A = x_A^*.$$}

In the notation, the arrow $\leftarrow$ and the double bar $||$ indicate conditioning by doing. In other words, $X \leftarrow x$ means that we intervene to force the value of the random variable $X$ to be $x$. In the intervention formula, conditioning by $x_A = x_A^*$ is equivalent to a substitution. Applying the intervention formula means that edges from parents of nodes in $A$ to nodes in $A$ are removed. This 'local surgery' yields a mutilated graph and the joint distribution factorizes along it. Note that we can also define an alternative formulation $P(X_v = x_v|X_A \leftarrow x_A^*) = P(X_{V \setminus A} = x_{V \setminus A}|X_A = x_A^*)$ as variables in $X_A$ have known values after the intervention.

A property of Do-Calculus is that if node $v$ has no parents then

$$P(X_v = x_v|X_v \leftarrow x_v^*) = P(X_v = x_v|X_v = x_v^*).$$

In this case, forcing a value results in same conditional probabilities as observing it. In other words, see and do probabilities are the same. Another property is the exogeneity:

$$P(X_v|X_{\pi(v)} \leftarrow x) = P(X_v|X_{\pi(v)} = x).$$
It means that forcing values to parents of a node \( v \) results in same conditional probability for \( v \) as observing those values in the parent variables \( X_{\pi(v)} \). Finally, we have the invariance property. For all \( v = 1, \ldots, n \) and \( S \subseteq V \) such that \( S \cap \{v, \pi(v)\} = \emptyset \), we have

\[
P(X_v | X_{\pi(v)} \leftarrow x, X_S \leftarrow s) = P(X_v | X_{\pi(v)} \leftarrow x).
\]

This means that once we control parents of \( v \), no further interventions will affect the probability of \( X_v \).

Here, we give an example of how an intervention affects a Bayesian network. The example is adapted from Pearl (2000) and Koski’s lecture slides from the Bayesian network course at KTH, year 2013. We are modeling slipperiness of a pavement and we assume that the slipperiness is a binary random variable denoted by \( X_5 \). So the pavement either is slippery or not. Slipperiness is affected by several environmental attributes which we model as random variables. These include the wetness of the pavement \( (X_4) \), sprinkler being possibly on \( (X_3) \), the possible rain \( (X_2) \), and the season \( (X_1) \) which is assumed to have four possible values. All the other random variables in this example are binary. Relations between these variables are described by the DAG of the Bayesian network. For example, given information about the wetness of the pavement, the slipperiness is independent of the random variables associated with the season, the rain and the sprinkler. The joint distribution of the random variables factorizes along the DAG:

\[
P(X_1 = x_1, X_2 = x_2, X_3 = x_3, X_4 = x_4, X_5 = x_5)
= P(X_1 = x_1)P(X_2 = x_2 | X_1 = x_1)P(X_3 = x_3 | X_1 = x_1) P(X_4 = x_4 | X_2 = x_2, X_3 = x_3) P(X_5 = x_5 | X_4 = x_4).
\]

The situation is illustrated in Figure 6.2.

Now we consider an intervention where the sprinkler is set on, i.e. we set \( X_3 \) equal to 1. We apply the intervention formula to get the do-conditioned probability:
6.2. DO-CALCULUS

\[ P(X_1 = x_1, X_2 = x_2, X_3 = x_3, X_4 = x_4, X_5 = x_5 | X_3 \leftarrow 1) \]
\[ = \frac{P(X_1 = x_1, X_2 = x_2, X_3 = 1, X_4 = x_4, X_5 = x_5)}{P(X_3 = 1 | X_1 = x_1)} \]
\[ = P(X_1 = x_1)P(X_2 = x_2 | X_1 = x_1)P(X_4 = x_4 | X_2 = x_2, X_3 = 1)P(X_5 = x_5 | X_4 = x_4) \].

The DAG in Figure 6.3 illustrates the situation. The edge from \( X_1 \) to \( X_3 \) has been removed. This is because after the intervention, i.e. setting the sprinkler on, the season has no effect on the sprinkler.

There are connections between Rothman’s sufficient-component cause model and Pearl’s Do-Calculus. As a part of probabilistic assessment of causal relations in the sufficient-component cause model, Do-Calculus could be used. We could describe a sufficient cause, i.e. a pie in the model, with a Bayesian network. Component causes would be parental nodes for the
outcome variable, i.e. state of an individual. Then do-conditioning would correspond to having an experiment where values of putative component causes are controlled. More complicated modeling would be needed to consider several sufficient causes. Also, component causes themselves could have causes behind them as represented by parental nodes. Generally, the idea of intervention in Do-Calculus is potentially useful for inference that relies on controlling values of some variables, i.e. conducting controlled experiments.
In this thesis, we have presented methods for sequence analysis and causal inference. In Article I, model class of Sparse Markov chains is defined and its properties are investigated. In Article II, a stochastic partition model of DNA sequences is considered so that the sequences are modeled as Markov chains of fixed order. Article III develops inference for SMC models further by introducing a recursive deterministic algorithm that uses Delaunay triangulation and Bayes factors. Finally, in Article IV causal inference is addressed and an algorithm for estimating interactions of multiple causes for a disease is given.

One underlying theme in this thesis is the use of the Bayesian approach for problems in machine learning. Especially, clustering is a task that we have addressed several times. A Bayesian approach to clustering is used in Articles I, II and III. Developing further the framework of Bayesian clustering could lead to improvements in the particular solutions we have presented in the articles. One possible direction is to use different prior distributions for partitions as now we have concentrated on the uniform prior and the Dirichlet process prior. For example, an uniform prior on the number of clusters could be used (Kohonen and Corander, 2014; Knorr-Held and Raßer, 2000; Quintana and Iglesias, 2003).

A second theme is the adaptation to the reality of large data sets. Our choice of learning algorithms reflects this. In Article I, a greedy stochastic algorithm is used instead of an MCMC approach that could provide consistent estimates (Marttinen et al., 2006; Corander et al., 2006, 2008). Even
the stochastic greedy algorithm proves to be relatively slow with very large data sets. In Article III, a heuristic deterministic recursive algorithm for the SMC models is presented. In future research, further approximations could be developed to deal with even larger data sets. In the articles, the use of point-estimation to obtain for example a MAP estimate has been a favoured method because estimating the posterior distribution would be computationally challenging. But for some future applications, estimating the posterior distribution with MCMC techniques could be valuable.

The data we have considered has been mostly sequence data. The use of Markovian models for sequence data is already a prominent tool for the scientist studying a variety of phenomena. As the data sets are growing, the need for sparse models will likely increase. The SMC model presented in Article I is a promising foundation for algorithms that process sequence data in a sparse manner. Article III provides a faster algorithm for learning an SMC model. Together these could be used in a variety of applications. For example, they could be combined with the clustering of Markov chains as presented in Article II.

Article IV presents a model for learning causal mechanisms from data. Although the Bayes’ rule is used as a part of the learning algorithm, we can say that overall the solution is not Bayesian. This shows a pragmatic attitude to probabilistic modeling. Generally, we use Bayesian methods because they work in practice but we are at the same time interested in non-Bayesian methodology as well. Also, using greedy and heuristic algorithms when facing big data sets can be considered as pragmatic.

There are some attractive areas of application for SMC models that however seem to be computationally too burdensome. For example, processing of natural language with any Markov model is challenging because of the memory required for software implementation. For the analysis of DNA data, there is an abundant literature of methods. Our approaches belong to the category of alignment free methods. For more on aligning DNA sequences, see Cheng et al. (2012).

For all the four articles, there has been a need to implement the algo-
rithms by developing software. For Article I, BAPS\(^1\) software was used as a part of the implementation. Based on Article II, a software package that is freely available was developed \(^2\). In future, software from Articles III and IV could also be made available for the public.

\(^1\)http://www.helsinki.fi/bsg/software/BAPS/
\(^2\)http://www.helsinki.fi/bsg/software/BACDNAS/


