# Bayesian Causal Inference 

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#### Abstract

In this thesis we address the problem of two-variable causal inference. This task refers to inferring an existing causal relation between two random variables (i.e. $X \rightarrow Y$ or $Y \rightarrow X$ ) from purely observational data. We begin by outlining a few basic definitions in the context of causal discovery, following the widely used do-Calculus $\overline{\mathrm{Pea} 00}$. We continue by briefly reviewing a number of state-of-the-art methods, including very recent ones such as CGNN Gou+17] and $K C D C$ MST18.

The main contribution is the introduction of a novel inference model where we assume a Bayesian hierarchical model, pursuing the strategy of Bayesian model selection. In our model the distribution of the cause variable is given by a Poisson lognormal distribution, which allows to explicitly regard discretization effects. We assume Fourier diagonal covariance operators, where the values on the diagonal are given by power spectra. In the most shallow model these power spectra and the noise variance are fixed hyperparameters. In a deeper inference model we replace the noise variance as a given prior by expanding the inference over the noise variance itself, assuming only a smooth spatial structure of the noise variance. Finally, we make a similar expansion for the power spectra, replacing fixed power spectra as hyperparameters by an inference over those, where again smoothness enforcing priors are assumed.

Based on our assumptions we describe an algorithmic forward model in order to produce synthetic causal data. These synthetic datasets are being used as benchmarks in order to compare our model to existing State-of-the-art models, namely LiNGAM [Hoy+09], ANMHSIC Moo+16], ANM-MML Ste+10], IGCI Dan+10] and CGNN Gou+17. We explore how well the above methods perform in case of high noise settings, strongly discretized data and very sparse data. Our model (BCM) shows to perform generally reliable in case of the synthetic data sets. While it is able to provide an accuracy close to the ANM methods in case of high noise and strongly discretized data, which deliver the best performance here, it is able to outperform other methods in case of very sparse ( 10 samples) synthetic data. We further test our model on the TCEP benchmark set, which is a widely used benchmark with real world data. Here our model can provide an accuracy comparable to state-of-the-art algorithms and is able to outperform other methods in a setting where only a small number of samples ( 20 samples) are available.


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## List of Symbols

| $\beta$ | a function, $\beta \in \mathbb{R}^{[0,1]}$ |
| :---: | :---: |
| $\delta_{z}$ | The Dirac delta distribution centered at $z \in \mathbb{R}$, i.e. $\delta_{z}=$ $\delta(\cdot-z)$ |
| $\eta$ | logarithmic noise variance, $\varsigma(x)^{2}=e^{\eta(x)}$ |
| $\epsilon$ | N -vector of noise samples $\left(\epsilon_{1}, \ldots, \epsilon_{N}\right)^{T}$ |
| $\lambda$ | $n_{\text {bins }}$-vector of Poisson means, $\boldsymbol{\lambda}=\left(\lambda_{1}, \ldots, \lambda_{n_{\text {bins }}}\right)^{T}$ |
| $\underset{c^{2}}{\Lambda(z)}$ | $N \times N$ matrix, with entries $\Lambda(z)_{i j}=\left(\mathcal{F}^{\dagger} \stackrel{e^{\tau_{f}} \delta_{z}}{ } \mathcal{F}\right)_{x_{i} x_{j}}$ variance of the noise in the i.i.d. noise case |
| $\sigma_{\beta / \eta / f}$ | hyper parameters controlling the strength of the smoothness or slope priors |
| $\tau_{\beta / f}$ | logarithmic power spectra, $P_{\beta / f}(q)=e^{\tau_{\beta / f}(q)}$ |
| B | covariance matrix of the distribution for $\beta$ |
| $d$ | tuple of the measurements $\boldsymbol{x}, \boldsymbol{y}$ |
| $\mathcal{E}$ | covariance matrix of the distribution for $\boldsymbol{\epsilon}$ |
| $f$ | a function representing the causal mechanism, $f \in \mathbb{R}^{[0,1]}$ |
| $F$ | covariance matrix of the distribution for $f$ |
| $\tilde{F}$ | $N \times N$ matrix, with entries $\tilde{F}_{i, j}=F\left(x_{i}, x_{j}\right)$ |
| $G$ | $N \times N$ matrix: $G=\left(\tilde{F}+\widehat{\left.e^{\eta(\boldsymbol{x}}\right)^{-1}}\right.$ |
| $\mathcal{F}$ | Fourier transformation operator |
| $\mathcal{H}(\cdot)$ | Information Hamiltonian, $\mathcal{H}(\cdot)=-\log \mathcal{P}(\cdot)$ |
| $k$ | $n_{\text {bins }}$-vector of counts ( $\left.k_{1}, \cdots, k_{n_{\text {bins }}}\right)^{T}$ |
| $\mathcal{N}(u \mid m, U)$ | Normal distribution in $u$, centered at $m$, with covariance U |
| $\mathcal{P}(\cdot)$ | A probability distribution, sometimes also being used directly as the corresponding probability density function |
| $\mathcal{P}(A \mid B)$ | the conditional probability of $A$ given $B$ |
| $\boldsymbol{x}$ | N -vector of measurements $\left(x_{1}, \cdots, x_{N}\right)^{T}$ |
| $y$ | N -vector of measurements $\left(y_{1}, \cdots, y_{N}\right)^{T}$ |
| $\int \mathcal{D}[u]$ | Notation for a path-integral over a function or field $u$ |
| $\sim$ | Notation for a random variable being distributed according to some probability distribution |
| $\hookleftarrow$ | Notation for a specific value of a random variable being drawn as a sample from some probability distribution |
| - $\mid$ | The determinant when applied to operators or matrices |
| , | Notation for a vector or field raised to a diagonal matrix or operator |

## 1. Introduction

### 1.1. Motivation and Significance of the Topic

Causal Inference regards the problem of drawing conclusions about how some entity we can observe does - or does not - influence or is being influenced by another entity. Having knowledge about such law-like causal relations enables us to predict what will happen ( $\hat{=}$ the effect) if we know how the circumstances ( $\hat{=}$ the cause) do change. For example, one can draw the conclusion that a street will be wet (the effect) whenever it rains (the cause). Knowing that it will rain, or indeed observing the rainfall itself, enables one to predict that the street will be wet. Less trivial examples can be found in the fields of epidemiology (identifying some bacteria as the cause of a sickness) or economics (knowing how taxes will influence the GDP of a country).

As PJS17 remark, the mathematical formulation of these topics has only recently been approached. Especially within the fields of data science and machine learning specific tasks from causal inference have been attracting much interest recently. HHH18 propose that causal inference stands as a third main task of data science besides description and prediction. Judea Pearl, best known for his Standard Reference Causality: Models, Reasoning and Inference, recently claimed that the task of causal inference will be the next "big problem" for Machine Learning Pea18. Such a specific problem is the two variable causal inference, also addressed as the cause-effect problem by PJS17. Given purely observational data from two random variables, $X$ and $Y$, which are directly causally related, the challenge is to infer the correct causal direction. In the example of rain and wetness of a street, this would mean, we are given two-dimensional observation samples corresponding to (Rainfall, Wetness of street). The samples itself could be ("it rains", "the street is wet"), ("it doesn't rain", "the street is wet"), ("it doesn't rain", "the street is dry"). We now have to conclude the true causal direction which is obviously Rainfall $\rightarrow$ Wetness of street.

Having only observational data means we can not intervene into the data, e.g. use a garden hose to see if the street gets wet when we simulate rainfall. In such a setting, inferring the true direction might seem to be a futile task. Indeed, inferring the true direction in the above example would be impossible if we only had observed the samples ("it rains", "the street is wet"), ("it doesn't rain", "the street is dry"). It was just the sample ("it doesn't rain", "the street is wet") that allowed us to discard the hypothetical direction Wetness of street $\rightarrow$ Rainfall. If we know that either one or the other direction have to be true, we can therefore conclude the true causal direction. Interestingly, this is an incorporation of a fundamental asymmetry between cause and effect which does always hold and can be exploited to tackle such an inference problem. Given two random variables, $X$ and $Y$ which are related
causally, $X \rightarrow Y$ (" $X$ causes $Y$ "), there exists a fundamental independence between the distribution of the cause $\mathcal{P}(X)$ and the mechanism which relates the cause $X$ to the effect $Y$. This independence however does not hold in the reverse direction. Most of the proposed methods for the inference of such a causal direction make use of this asymmetry in some way, either by considering the independence directly [Dan+10], Moo+16], or by taking into account the algorithmic complexity for the description of the factorization $\mathcal{P}(X) \mathcal{P}(Y \mid X)$ and comparing it to the complexity of the reverse factorization $\mathcal{P}(Y) \mathcal{P}(X \mid Y)$.

### 1.2. Structure of the Work

The rest of the thesis will be structured as following. In Chapter 2 we will first outline and specify our problem setting. We will attempt to define a self-contained framework of definitions for the necessary models of causality, following the do-Calculus introduced by $\mathrm{Pea00}$. We also will review existing methods here, namely Additive Noise Models, Information Geometric Causal Inference and Learning Methods.

Chapter 3 will describe our inference model which is based on a hierarchical Bayesian model. In Section 3.2 we will introduce a first, shallow inference model. Here we assume covariance operators, determining distributions for the causal mechanism and the cause variable distribution itself, as well as the noise variance, to be given.

In Section 3.3 we will relax the fixed assumptions by allowing the noise variance to be determined by a field itself which only is governed by a prior distribution controlling the slope of the field. This will be expanded in 3.4 by allowing for arbitrary power spectra which are assumed to be random variables distributed by smoothness enforcing priors.

In Chapter 4 we will accompany the theoretical framework with experimental results. As the computational implementation of the deeper models $(3.3,3.4)$ showed to be problematic, we will limit the considerations here to an implementation of the shallow inference model proposed in 3.2. To that end we begin by outlining a "forward model" which allows to sample causally related data in 4.1. We describe a specific algorithm for the inference model in 4.2, which is then tested on various benchmark data (4.4). To provide a reference, the performance is compared to state-of-the-art methods introduced in 2.2 .

We conclude by assessing that our model generally can show competitive classification accuracy and propose possibilities to further advance the model.

### 1.3. Related Work

## 2. Problem Setting and Related Work

### 2.1. Problem Setting

We begin by briefly defining key concepts in causal inference using the do-Calculus, introduced by Pea00]. This outline attempts to be as short as possible as necessary for the present thesis, however still self contained.

Definition 1. A causal structure of some random variables $\mathcal{V}=\left\{X_{1}, \ldots, X_{n}\right\}$ is a directed acyclic graph (DAG) $\mathcal{D}=(\mathcal{V}, \mathcal{E})$ with the elements of $\mathcal{V}$ as nodes and the edges $\mathcal{E}$ representing functional relations between the variables.

Returning to the example of rain causing a street to be wet from the beginning, one could describe this situation by a causal structure with two vertices ( $R$ for rainfall and $W$ for wetness of the street) and one directed edge ( $R, E$ ). ${ }^{\text {® }}$ To make this more illustrating, we expand the model by considering a sprinkler next to the street, with a state $S$ either being on or off, as another possible cause for the street to be wet, and we consider the cloudiness of the sky $(C)$ as a "reason" for the rainfall. The causal structure is therefore given by the DAG $\mathcal{D}=(\mathcal{V}, \mathcal{E})$ with $\mathcal{V}=\{C, R, S, W\}$ and $\mathcal{E}=\{(C, R),(R, W),(S, W)\}$. The full graph structure is depicted in Fig. 2.1.


Figure 2.1.: An example for a DAG representing a causal structure, the nodes are random variables and the arrows give the directed edges, representing functional relations.

Definition 2. A pair $\mathcal{M}=\left(\mathcal{D}, \Theta_{\mathcal{D}}\right)$ of a graph $\mathcal{D}=(\mathcal{V}, \mathcal{E})$ and its parameters $\Theta_{\mathcal{D}}$ defines a causal model. The parameters $\Theta_{\mathcal{D}}$ assign to each $X_{i} \in \mathcal{V}$ an equation of the form

$$
\begin{equation*}
x_{i}=f\left(p a_{i}, u_{i}\right) \tag{2.1}
\end{equation*}
$$

[^0]and a probability distribution $\mathcal{P}\left(U_{i}\right)$ where $P a_{i}$ denote the parents of $X_{i}$ w.r.t. to the graph $\mathcal{D}$ and $U_{i}$ are unobserved disturbances having influences on $X_{i} . x_{i}, p a_{i}, u_{i}$ denote realizations of $X_{i}, P a_{i}$ and $U_{i}$ respectively.

In the above example (Fig. 2.1), the vertices $C, S$ do not have any (observed) parental nodes, so their values are only determined by unobserved noise $\left(c=f_{c}\left(u_{c}\right)\right.$ and $\left.s=f_{s}\left(u_{s}\right)\right)$. We think of the rainfall as an effect of the sky's cloudiness, so the corresponding equation would be $r=f_{r}\left(c, u_{r}\right)$, where we allow for some unobserved nuisance. The wetness of the street finally does have two separate observed causes - the rainfall and the sprinkler. This corresponds to an equation of the form $w=f_{w}\left(r, s, u_{w}\right)$, again allowing for some unobserved "noise" influencing the measured state of the street wetness. Now we can define an intervention in the model which represents an external manipulation of some variable:

Definition 3. Given some causal model $\mathcal{M}$, with a random variable $X_{i} \in \mathcal{V}$ the atomic intervention $\boldsymbol{d o}\left(X_{i}=x_{i}\right)$ is defined by setting the value of $X_{i}$ to $x_{i}$, removing the equation regarding $x_{i}=f\left(p a_{i}, u_{i}\right)$ from the parameters of the model and substituting $X=x_{i}$ in all other equations.

Assuming for a moment here we had the power to make it rain, we could set the variable $R=$ rainfall in our example. This would remove the equation $r=f_{r}\left(c, u_{r}\right)$ from the parameters and set $w=\left(R=\right.$ rainfall, $\left.s, u_{w}\right)$

Using this do-formalism we can provide a definition for one variable being the cause of another one:

Definition 4. Given two random variables $X, Y$ with a joint probability distribution $\mathcal{P}(X, Y)$ and a corresponding conditional distribution $\mathcal{P}(Y \mid X)$ we say $X$ causes $Y$ (denoted by $X \rightarrow Y$ ) iff $\mathcal{P}(y \mid \operatorname{do}(x)) \neq \mathcal{P}\left(y \mid \operatorname{do}\left(x^{\prime}\right)\right)$ for some $x, x^{\prime}$ being realizations of $X$ and $y$ being a realization of $Y$

For a last time returning to the above example, the probability for the street being wet, when we force rainfall, $\mathcal{P}(w \mid \mathrm{do}(R=$ it rains $)$ will probably be much higher than the corresponding one for manipulating the system such that it doesn't rain, $\mathcal{P}(w \mid \operatorname{do}(R=\mathbf{i t}$ doesn't rain $)$. In contrast to this, watering the street with a garden hose and thus making it wet will most likely not influence the probability of rain, $\mathcal{P}(r \mid \operatorname{do}(w=$ the street is wet $)=\mathcal{P}(r \mid \operatorname{do}(w=$ the street is dry)

Fig. 2.2 shows some possibilities in which way two variables can be causally related. In case (c) there is no causal relation at all. Note that case (d) is not consistent with our definition of causality as it cannot be modeled with a DAG. Case (e) corresponds to a confounding variable, i.e. a variable that influences both of the other ones. A standard example in the statistical literature would be the human birth rate and the stork population, where one can find a significant correlation Mat00. As it is however unlikely that one of these variables directly causes the other one, there might be some confounding variable as industrialization or environmental health (and therefore standards of living) which is causing the other ones [ST05].

Figure 2.2.: Models for causal relations in the 2 variable case, reproduced from Moo +16

(a) $X \rightarrow Y$
$\mathcal{P}(Y) \neq \mathcal{P}(Y \mid \operatorname{do}(X=x))=\mathcal{P}(Y \mid X=x)$

$$
\mathcal{P}(X)=\mathcal{P}(X \mid \operatorname{do}(Y=y)) \neq \mathcal{P}(X \mid Y=y)
$$


(c) no causal relation between $X$ and $Y$

$$
\begin{aligned}
& \mathcal{P}(Y)=\mathcal{P}(Y \mid \operatorname{do}(X=x))=\mathcal{P}(Y \mid X=x) \\
& \mathcal{P}(X)=\mathcal{P}(X \mid \operatorname{do}(Y=y))=\mathcal{P}(X \mid Y=y)
\end{aligned}
$$


(e) (Hidden) confounder, $Z \rightarrow X$ and $Z \rightarrow Y$

$$
\begin{aligned}
& \mathcal{P}(Y)=\mathcal{P}(Y \mid \operatorname{do}(X=x)) \neq \mathcal{P}(Y \mid X=x) \\
& \mathcal{P}(X)=\mathcal{P}(X \mid \operatorname{do}(Y=y)) \neq \mathcal{P}(X \mid Y=y)
\end{aligned}
$$


(b) $Y \rightarrow X$
$\mathcal{P}(Y)=\mathcal{P}(Y \mid \operatorname{do}(X=x)) \neq \mathcal{P}(Y \mid X=x)$

$$
\mathcal{P}(X) \neq \mathcal{P}(X \mid \operatorname{do}(Y=y))=\mathcal{P}(X \mid Y=y)
$$


(d) $X \rightarrow Y$ and $Y \rightarrow X$ $\mathcal{P}(Y) \neq \mathcal{P}(Y \mid \operatorname{do}(X=x)) \neq \mathcal{P}(Y \mid X=x)$ $\mathcal{P}(X) \neq \mathcal{P}(X \mid \operatorname{do}(Y=y)) \neq \mathcal{P}(X \mid Y=y)$

(f) V-collider, $X \rightarrow S$ and $Y \rightarrow S$
$\mathcal{P}(Y \mid S=s) \neq \mathcal{P}(Y \mid \operatorname{do}(X=x), S=s)=\mathcal{P}(Y \mid X=x, S=s)$
$\mathcal{P}(X \mid S=s) \neq \mathcal{P}(X \mid \operatorname{do}(Y=y), S=s)=\mathcal{P}(X \mid Y=y, S=s)$

Case (f) is often referred to as a V-Collider [Spi16] and illustrates the problem of selection bias. A car mechanic, only working with cars which do not start $(S=0)$, observes whether the start engine is broken $(X=0)$ and if the battery is empty $(Y=0)$. As in most cases one of these might be the case, but not the other one as well, he might draw the conclusion that $X$ and $Y$ are causally related (if $X=1$, then usually $Y=0$ and vice versa if $Y=1$ then $X=0$ ), if he does not consider his conditioning on the selection bias $S=0$ Moo+16]. I.e. the mechanic considers only $\mathcal{P}(X \mid Y), \mathcal{P}(Y \mid X)$ instead of $\mathcal{P}(X \mid Y, S), \mathcal{P}(Y \mid X, S)$ Not taking such a selection bias into account can thus lead to conclude a non-existing causal relation.

In this work we however only consider case (a) and (b), i.e. $X$ being a cause of $Y$ and $Y$ being a cause of $X$ and deciding which is the true one. Having access to the distributions $\mathcal{P}((X) \mid \operatorname{do}(Y))$ and $\mathcal{P}(Y \mid \operatorname{do}(X))$ would render this an easy task, as one could simply check via the definition 4 which direction holds. However, such distributions are usually not available, as manipulation of the systems of interest is often not possible. Instead we assume to have continuous observations in form of samples $\left(x_{i}, y_{i}\right)$ being drawn from $\mathcal{P}(X, Y)$ and want to know which of the cases, (a) or (b) in Fig. 2.2 holds for the underlying process. We state the problem as following:

```
Problem 1. Prediction of causal direction for two variables
Input: A finite number of sample data \(\boldsymbol{d} \equiv(\boldsymbol{x}, \boldsymbol{y})\), where \(\boldsymbol{x}=\left(x_{1}, \ldots, x_{N}\right), \boldsymbol{y}=\left(y_{1}, \ldots, y_{N}\right)\)
Output: A predicted causal direction \(\mathcal{D}_{X \rightarrow Y} \in\{-1,1\}\) where -1 represents the prediction
" \(Y \rightarrow X\) " and 1 represents " \(X \rightarrow Y\) "
```


### 2.2. Related Work

Approaches to tasks in causal inference from purely observational data are often divided into three groups ([SZ16; MST18]), namely constraint-based, score-based and asymmetry-based methods. Sometimes this categorization is extended by considering learning methods as a fourth, separate group. Two of these categories, constraint-based methods and score-based methods are basically searching for the true DAG representing some structure and rely on tests of conditional independence using conditioning on external variables. As such are not available in the two-variable case, those models are of little interest for the present task.

A third category exploits an inherent asymmetry between cause and effect. This asymmetry can be framed in different terms. One way is to use the concept of algorithmic complexity - given a true direction $X \rightarrow Y$, the factorization of the joint probability into $\mathcal{P}(X, Y)=\mathcal{P}(X) \mathcal{P}(Y \mid X)$ will be less complex than the reverse factorization $\mathcal{P}(Y) \mathcal{P}(X \mid Y)$ This approach is often used by Additive Noise Models (ANMs). Another way is to state that the mechanism relating cause and effect should be independent of the cause Dan+10]. This formulation is employed by the concept of Information Geometric Causal Inference.

We will also consider Causal Generative Neural Networks (CGNN) as an example of learning methods. For the following, let $\mathcal{X}, \mathcal{Y}$ be some measurable spaces. We consider some random variables $X$ and $Y$ on $\mathcal{X}$ and $\mathcal{Y}$.

### 2.2.1. Additive Noise Models - Principles

A large family of inference models assume additive noise, i.e. in the case $X \rightarrow Y, Y$ is determined by some function $f$, mapping $X$ to $Y$, and some collective noise variable $E_{Y}$, i.e. $Y=f(X)+E_{Y}$, where $X$ is independent of $E_{Y}$. Hypothetically, the same can be done in the backwards direction: $X=g(Y)+E_{X}$. Moo +16 show that the resulting joint distribution $\mathcal{P}(X, Y)$ of such an ANM is either induced by the forward or the backward model, but generally not by both. If this is the case the model is said to be identifiable.

### 2.2.2. LiNGAM - an ICA-based Approach

An early adoption of this principle, $\operatorname{LiNGAM}$ (Linear Non-Gaussian Additive Noise Model, [Shi+06]) models the structure with linear functions, i.e. writing $x=c_{x} y+e_{x}$ or $y=c_{y} x+e_{y}$. This can be modeled using a vector calculus

$$
\begin{equation*}
\binom{x}{y}=C\binom{x}{y}+\binom{e_{x}}{e_{y}} \tag{2.2}
\end{equation*}
$$

As the restriction in the given problem setting allows only one direction of $X \rightarrow Y, Y \rightarrow X$ to be true, one of $c_{x}, c_{y}$ can be set to zero, meaning $C$ is a triangular matrix. Introducing the mixing matrix $A \equiv(\mathbb{1}-C)^{-1}$, the above relation can be re-written as

$$
\begin{equation*}
\binom{x}{y}=A\binom{e_{x}}{e_{y}} \tag{2.3}
\end{equation*}
$$

The authors assume non-Gaussian error terms $e_{x}, e_{y}$ which makes it possible to employ the technique of independent component analysis (ICA) to estimate the mixing matrix $A$ which yields an estimation of the component matrix $C=\mathbb{1}-A^{-1}$. Finally a strictly lower triangular permutation is sought, the permutation matrix gives the causal ordering of the variables.

The model is proposed for a multi-variable case, the two-variable case is actually just a special case here. Even though we do not make the two main assumptions (namely linear relations and non-Gaussianity of the errors) in the course of specifying our model, it is mentioned at this place as LiNGAM has become a standard reference for causal inference benchmarks Ste+10; MST18.

### 2.2.3. Additive Noise Models with Regression and Residuum Independence Scoring

More recent approaches usually deal with the possibility of non-linear causal mechanisms. As mentioned above, there is a fundamental asymmetry between cause and effect which incorporates in the fact that the distribution $\mathcal{P}(X)$ of the cause is independent from the true causal mechanism $f$ itself. This functional mechanism is represented by the conditional distribution $\mathcal{P}(Y \mid X)$. A way to exploit this is to make a regression $\hat{f}$ for $f$ and calculate the
residual $\hat{f}(X)-Y$. If $\hat{f}$ would be an exact regression of $f$ and $X \rightarrow Y$, this residual would now be fully independent of $X$.

However the regression is only somewhat precise (Gaussian Process Regression is usually employed here $[$ Hoy $+09 ;$ Moo+16] ) and not the whole distribution $\mathcal{P}(X)$ but only a finite number of samples $\left(x_{1}, \ldots, x_{N}\right)$ is given. Therefore one cannot expect to measure complete independence between $\left(x_{1}, \ldots, x_{N}\right)$ and $\left(\hat{f}\left(x_{1}\right)-y_{1}, \ldots, \hat{f}\left(x_{N}\right)-y_{N}\right)$ but only a higher independence score in the true direction. A number of proposed methods use the (empirical) Hilbert Schmidt Independence Criterion (HSIC) to estimate the the independence. This measure is defined using the formulation of Reproducing Kernel Hilbert Spaces (RKHS), in which probability distributions can be bijectively embedded as elements.

An explicit introduction of this framework would not be in scope within this thesis, however we want to give a consistent, self-contained overview of the key concept used by recent publications in causal inference. The following outline is based on Gre+05 and Gre+07.

For a measurable space $\mathcal{X}$, the RKHS $\mathcal{H}$ is a Hilbert space of functions $f: \mathcal{X} \rightarrow \mathbb{R}$ with an inner product $\langle\cdot, \cdot\rangle_{\mathcal{H}}$, in which the evaluation functionals $\delta_{x}: f \mapsto f(x)$ are continuous for all $x \in \mathcal{X}$. As guaranteed by the Riesz representation theorem theorem, one can always represent such an evaluation at some $x \in \mathcal{X}$ by taking the inner product with a unique element of the Hilbert space:

$$
\begin{equation*}
f(x)=\delta_{x}(f)=\left\langle k_{x}, f\right\rangle_{\mathcal{H}} \tag{2.4}
\end{equation*}
$$

The function $k_{x}$ therefore defines a kernel via $k_{x}(y)=\left\langle k_{y}, k_{x}\right\rangle_{\mathcal{H}} \equiv k(x, y)$ and thus fully specifies the RKHS, allowing us to write $\mathcal{H}_{k}$ A probability distributions $\mathcal{P}$ over $\mathcal{X}$ can be embedded into such a RKHS $\mathcal{H}_{\mathcal{X}}$ by the mean embedding:

Definition 5. For a kernel $k: \mathcal{X} \times \mathcal{X} \rightarrow \mathbb{R}$, a random variable $X \sim \mathcal{P}$ in $\mathcal{X}$, the mean embedding $\mu_{k}(\mathcal{P})$ is defined by:

$$
\begin{equation*}
\mu_{k}(\mathcal{P}) \equiv \mathbb{E}_{X \sim \mathcal{P}}[k(\cdot, X)] \tag{2.5}
\end{equation*}
$$

given samples $\left(x_{1}, \ldots, x_{N}\right)$ from $X$ an estimation of the mean embedding is given by

$$
\begin{equation*}
\hat{\mu}_{k}(\mathcal{P}) \equiv \frac{1}{N} \sum_{i=1}^{N} k\left(\cdot, x_{i}\right) \tag{2.6}
\end{equation*}
$$

An inner product with the mean embedding is therefore given by Eq. 2.4 via

$$
\begin{equation*}
\left\langle f, \mu_{k}(\mathcal{P})\right\rangle_{\mathcal{H}_{k}} \equiv \mathbb{E}_{X \sim \mathcal{P}}[f(X)] \tag{2.7}
\end{equation*}
$$

Now one can define the HSIC Gre+05

Definition 6. Given kernels $k: \mathcal{X} \times \mathcal{X} \rightarrow \mathbb{R}$ and $y: \mathcal{Y} \times \mathcal{Y} \rightarrow \mathbb{R}$, the HSIC is defined as

$$
\begin{align*}
& \operatorname{HSIC}(X, Y)_{k, l} \equiv \| \mu_{k \otimes l}(\mathcal{P}(X, Y))-\mu_{k \otimes l}\left(\mathcal{P}(X) \mathcal{P}(Y) \|_{\mathcal{H}_{k \otimes l}}=\right. \\
&=\mathbb{E}_{X, X^{\prime} \sim \mathcal{P}(X)}\left[k\left(X, X^{\prime}\right) l\left(Y, Y^{\prime}\right)\right. \\
& Y Y^{\prime} \sim \mathcal{P}(Y) \\
&+\mathbb{E}_{X, X^{\prime} \sim \mathcal{P}(X)}\left[k\left(X, X^{\prime}\right)\right] \mathbb{E}_{Y, Y^{\prime} \sim \mathcal{P}(Y)}\left[l\left(Y, Y^{\prime}\right)\right]  \tag{2.8}\\
&-2 \mathbb{E}_{X \sim \mathcal{P}(X)}\left[\mathbb{E}_{X^{\prime} \sim \mathcal{P}(X)}\left[k\left(X, X^{\prime}\right)\right] \mathbb{E}_{Y^{\prime} \sim \mathcal{P}(Y)}\left[l\left(Y, Y^{\prime}\right)\right]\right] \\
& Y \sim \mathcal{P}(Y)
\end{align*}
$$

An estimation based on finite sample data $\boldsymbol{x}=\left(x_{1}, \ldots, x_{N}\right), \boldsymbol{y}=\left(y_{1}, \ldots, y_{N}\right)$, is given by $\widehat{\operatorname{HSIC}}(\boldsymbol{x}, \boldsymbol{y})_{k, l}=(N-1)^{-2} \operatorname{tr}(K H L H)$, with the Gram matrices $K_{i j}=k\left(x_{i}, x_{j}\right), L_{i j}=$ $l\left(x_{i}, x_{j}\right)$ and the centering matrix $H_{i j}=\delta_{i j}-N^{-1}$.

This method, in the following called $A N M-H S I C$, performed strongly in recent benchmarks ( Moo+16, Gou+17; MST18]). The authors Moo+16 use squared exponential (also called Gaussian) kernels here, i.e. $k(x, y)=e^{-\gamma(x-y)^{2}}$, where they allow the bandwidth $\gamma$ to be estimated from the data itself.

### 2.2.4. Empirical Bayes - Additive Noise with MML Scoring

Other models, also based on the additive noise model, i.e. the assumption $Y=f(X)+E$, use Bayesian model selection. Here one compares the probability $\mathcal{P}(X \rightarrow Y \mid \boldsymbol{d})$ to the probability of the competing direction $\mathcal{P}(Y \rightarrow X \mid \boldsymbol{d})$, where again $\boldsymbol{d}$ denotes the observed samples, $\boldsymbol{d}=(\boldsymbol{x}, \boldsymbol{y})$. The ratio of these model probabilities

$$
\begin{equation*}
\mathcal{O}_{X \rightarrow Y}=\frac{\mathcal{P}(X \rightarrow Y \mid \boldsymbol{d})}{\mathcal{P}(Y \rightarrow X \mid \boldsymbol{d})} \tag{2.9}
\end{equation*}
$$

is often referred to as the odds ratio or Bayes factor (see e.g. BS09).
Using Bayes Theorem

$$
\begin{equation*}
\mathcal{P}(X \rightarrow Y \mid \boldsymbol{d})=\frac{\mathcal{P}(X \rightarrow Y) \mathcal{P}(\boldsymbol{d} \mid X \rightarrow Y)}{\mathcal{P}(\boldsymbol{d})} \tag{2.10}
\end{equation*}
$$

and the fact that the prior probabilities of the competing models should be equal $(\mathcal{P}(X \rightarrow$ $Y)=\mathcal{P}(Y \rightarrow X))$ we can express 2.9 in terms of the marginal likelihoods:

$$
\begin{equation*}
\mathcal{O}_{X \rightarrow Y}=\frac{\mathcal{P}(\boldsymbol{d} \mid X \rightarrow Y)}{\mathcal{P}(\boldsymbol{d} \mid Y \rightarrow X)} \tag{2.11}
\end{equation*}
$$

Such a model is employed by Ste+10 which assume a Gaussian mixture model for the distribution of the cause $\mathcal{P}(X)$ and a Gaussian Process with a squared exponential kernel for the causal mechanism. Here, the authors include a regression step for the function $f$ representing the causal mechanism. Numerical quantities for the resulting terms are given by an expansion based on the concept of Minimum Message Length (MML). We will therefore refer to this approach as $A N M-M M L$.

### 2.2.5. Kernel Deviance Measures

In a very recent (April 2018) publication, MST18 introduced the method of Kernel Conditional Deviance. Here again the asymmetry in the algorithmic complexity between the factorizations of the joint probability $\mathcal{P}(X, Y), \mathcal{P}(X) \mathcal{P}(Y \mid X)$ and $\mathcal{P}(Y) \mathcal{P}(X \mid Y)$ is considered. The authors reason that in case the true causal direction is $X \rightarrow Y$, it holds that

$$
\begin{equation*}
K\left(\mathcal{P}\left(Y \mid x_{i}\right)\right)=K\left(\mathcal{P}\left(Y \mid x_{j}\right)\right) \quad \forall i, j \tag{2.12}
\end{equation*}
$$

which however is not true in the other direction:

$$
\begin{equation*}
K\left(\mathcal{P}\left(X \mid y_{i}\right)\right) \neq K\left(\mathcal{P}\left(X \mid y_{j}\right)\right) \quad \forall i, j \tag{2.13}
\end{equation*}
$$

Above, $K(\mathcal{P})$ denotes the Kolmogorov complexity of $\mathcal{P}$ which is, loosely speaking, the length of a program that encodes the distribution $\mathcal{P}[\mathrm{GV} 08]$. As this is not further specified, the Kolmogorov complexity itself is uncomputable and to be understood in a conceptional way. Based on this thought and using the variance in the of the conditional mean embedding of distributions in the RKHS, one derives at the estimator

$$
\begin{equation*}
S_{X \rightarrow Y}^{\mathrm{KCDC}}=\frac{1}{N} \sum_{i=1}^{N}\left(\left\|\mu_{Y \mid X=x_{i}}\right\|_{\mathcal{H}_{Y}}-\frac{1}{N} \sum_{j=1}^{N}\left\|\mu_{Y \mid X=x_{j}}\right\|_{\mathcal{H}_{\mathcal{Y}}}\right) \tag{2.14}
\end{equation*}
$$

and equivalently in the other direction $S_{Y \rightarrow X}^{\mathrm{KCDC}}$, with roles of $X$ and $Y$, resp. $x$ and $y$ switched. The direction predicted is the one with the lower deviance, i.e. $X \rightarrow Y$ if $S_{X \rightarrow Y}^{\mathrm{KCDC}}<S_{Y \rightarrow X}^{\mathrm{KCDC}}$, $Y \rightarrow X$ otherwise. The authors measure the performance of their algorithm in a experimental setup, describing mostly perfect predictions in case of synthetic data and very good ( $\approx 74 \%$ ) accuracy for the real world TCEP-benchmark.

### 2.2.6. Information Geometric Causal Inference

The concept that the distribution of the cause variable $\mathcal{P}(X)$ and the causal mechanism relating cause and effect, represented by the conditional distribution $\mathcal{P}(Y \mid X)$, represent independent mechanisms of nature is also the foundation of the approach of Information Geometric Causal Inference (IGCI). This approach has been introduced by Dan+10], where, instead of approximating algorithmic complexity, the orthogonality of independent distributions in information space is exploited. The authors show that their approach even works in a deterministic, noise-free scenario. In such a scenario, one has $X=f(Y)$ and because of the bijectivity one also can state $Y=f^{-1}(X)$. The concept of IGCI is the thought that given independence between $\mathcal{P}(X)$ and $f$, the covariance of $\mathcal{P}(X)$ and $\left|\log \left(f^{\prime}\right)\right|$ being considered as random variables, should vanish. Here the covariance w.r.t. some reference distribution is considered, in the simplest case, the uniform distribution on $[0,1]$. The authors arrive at the score

$$
\begin{equation*}
C_{X \rightarrow Y}^{\mathrm{IGCI}} \equiv \int \mathrm{~d} x \log \left(\left|f^{\prime}(x)\right|\right) \mathcal{P}(x)-\int \mathrm{d} y \log \left(\left|g^{\prime}(y)\right|\right) \mathcal{P}(y)=-C_{Y \rightarrow X}^{\mathrm{IGCI}} \tag{2.15}
\end{equation*}
$$

and infer $X \rightarrow Y$ whenever $C_{X \rightarrow Y}^{\mathrm{IGCI}}>0$ and $Y \rightarrow X$ otherwise. The authors show that Eq. 2.15 is is equivalent to the difference of Shannon entropies, i.e.

$$
\begin{equation*}
C_{X \rightarrow Y}^{\mathrm{IGCI}}=S(\mathcal{P}(X))-S(\mathcal{P}(Y)) \tag{2.16}
\end{equation*}
$$

with the Shannon entropy

$$
\begin{equation*}
S(\mathcal{U}) \equiv-\int \mathrm{d} x \mathcal{U}(x) \log \mathcal{U}(x) \tag{2.17}
\end{equation*}
$$

which they estimate on finite sample data $\boldsymbol{x}, \boldsymbol{y} \in[0,1]^{N}$ via:

$$
\begin{equation*}
\hat{S}(\boldsymbol{x}) \equiv \psi(N)-\psi(1)+\frac{1}{N-1} \sum_{i=1}^{N-1} \log \left|x_{i+1}-x_{i}\right| \tag{2.18}
\end{equation*}
$$

where $\psi(x)=\frac{d}{d x} \log \Gamma(x)$, the Digamma function. In 2.18, the samples are assumed to be in non-decreasing order, i.e. $x_{i} \leq x_{i+1}$. Further, the convention $\log (0)=0$ is assumed, so that repeated samples with $x_{i+1}-x_{i}=0$ do not contribute to $\hat{S}(\boldsymbol{x})$.

As a main drawback we consider the restriction to noiseless (or almost noiseless, the authors show that their work can be extended to models with small noise) case. Also the independence with respect to certain reference distributions is essential, as pointed out by the authors themselves. Furthermore, as [SZ16] remarks, the introduced method also relies on $\mathcal{P}(X)$ and $\left|\log \left(f^{\prime}\right)\right|$ being complex enough that they can be assessed for empirical results.

### 2.2.7. Learning Methods - CGNN

A recent publication [Gou+17] introduced a neural network-based approach for causal inference. The authors use neural networks with one hidden layer and a $R e L U$ activation function (see e.g. Bis06] for a detailed introduction on neural networks), taking samples from one variable, e.g. $X$, as the input and fitting the output to the other variable $Y$. The loss function is the (empirical) Maximum Mean Discrepancy, introduced by Gre+07:

Definition 7. Given a kernel $k: \mathcal{X} \times \mathcal{X} \rightarrow \mathbb{R}$ and random variables $X, Y$ on $\mathcal{X}$, the MMD can be defined as

$$
\begin{equation*}
\operatorname{MMD}_{k}(X, Y)=\left\|\mathbb{E}_{X} k(\cdot, X)-\mathbb{E}_{Y} k(\cdot, Y)\right\|_{\mathcal{H}_{k}} \tag{2.19}
\end{equation*}
$$

and can be estimated from samples $\boldsymbol{x}=\left(x_{1}, \ldots, x_{N}\right) \hookleftarrow X, \boldsymbol{y}=\left(y_{1}, \ldots, y_{N^{\prime}}\right) \hookleftarrow Y$ as the empirical MMD, given by:

$$
\begin{equation*}
\widehat{\operatorname{MMD}}_{k}(\boldsymbol{x}, \boldsymbol{y})=\frac{1}{N^{2}} \sum_{i j}^{N} k\left(x_{i}, x_{j}\right)+\frac{1}{N^{2}} \sum_{i j}^{N^{\prime}} k\left(y_{i}, y_{j}\right)-\frac{2}{N N^{\prime}} \sum_{i}^{N} \sum_{j}^{N^{\prime}} k\left(x_{i}, y_{j}\right) \tag{2.20}
\end{equation*}
$$

The authors use a Gaussian kernel where the bandwidth $\gamma$ is a hyperparameter to be set. After a training phase, in which the neural network is being tuned to predict the distribution
for $Y$ given the samples from $X$, the empirical MMD between the given samples $\boldsymbol{y}$ and the samples $\hat{\boldsymbol{y}}(\boldsymbol{x})$ predicted by neural network is being measured. The same is done for the reverse direction, $Y \rightarrow X$. The direction with the smaller MMD is then the preferred one. A known advantage of this method is that neural networks are universal approximators, meaning that basically every function can be approximated. The authors find that their method performs quite well, outperforming all other methods in the real-world data $T C E P$-benchmark.

## 3. Bayesian Inference Model

### 3.1. Our Contribution and the Formalism of IFT

Our contribution incorporates the concept of Bayesian model selection. As already briefly outlined in 2.2.4 this concept compares two competing models, in our case $X \rightarrow Y$ and $Y \rightarrow X$, and asks for the ratio of the marginalized likelihoods,

$$
\mathcal{O}_{X \rightarrow Y}=\frac{\mathcal{P}(\boldsymbol{d} \mid X \rightarrow Y, M)}{\mathcal{P}(\boldsymbol{d} \mid Y \rightarrow X, M)}
$$

Where $M$ denotes the hyperparameters which are assumed to be the same for both models.
In the setting of the present causal inference problem, such an approach has already been used by $[\mathrm{Ste}+10]$. In contrast to the above publication we will use the formalism of information field theory (IFT), introduced by EFK09.

IFT considers signal fields $s$ which reflect a physical state $\psi, s=s[\psi]$ and follows some probability $s \hookleftarrow \mathcal{P}(s)$. Such signal fields usually have infinite degrees of freedom, this makes them an adequate choice to model our distribution of the cause variable and the function relating cause and effect.

Throughout the following we will consider $X \rightarrow Y$ as the true underlying direction which we derive our formalism on. The derivation for $Y \rightarrow X$ will follow analogously by switching the variables.

### 3.2. A Shallow Inference Model

We will begin with deriving in 3.2 .1 the distribution of the cause variable, $\mathcal{P}(X \mid X \rightarrow Y, M)$ where $M$ defines a set of assumptions and hyperparameters we impose on the model and are yet to be specified. In 3.2 .2 we continue by considering the conditional distribution $\mathcal{P}(Y \mid X, X \rightarrow Y, M)$. Combining those results, we compute then the full Bayes factor in 3.2.3.


Figure 3.1.: Overview over the most shallow Bayesian hierarchical model considered, for the case $X \rightarrow Y$

### 3.2.1. Distribution of the Cause Variable

## Basic Considerations

Without imposing any constraints, we reduce our problem to the interval $[0,1]$ by assuming that $\mathcal{X}=\mathcal{Y}=[0,1]$. This can always be ensured by rescaling the data. Now we make the assumption that in principle, the cause variable $X$ follows a lognormal distribution.

$$
\begin{equation*}
\mathcal{P}(x \mid \beta) \propto e^{\beta(x)} \tag{3.1}
\end{equation*}
$$

with $\beta \in \mathbb{R}^{[0,1]}$, being some signal field which follows a zero-centered normal distribution, $\beta \sim \mathcal{N}(\beta \mid 0, B)$.
Here we write $B$ for the covariance operator $\mathbb{E}_{\beta \sim \mathcal{P}(\beta)}\left[\beta\left(x_{0}\right) \beta\left(x_{1}\right)\right]=B\left(x_{0}, x_{1}\right)$. We note that this is equivalent to the definition of $\beta$ as a Gaussian Process which would be $\beta(x)=$ $\mathcal{G} \mathcal{P}\left(0, B\left(x, x^{\prime}\right)\right)$ using the notation of RW06.

We postulate statistical homogeneity ${ }^{1}$ ). for the covariance, that is

$$
\begin{array}{r}
\mathbb{E}_{\beta \sim \mathcal{P}(\beta)}[\beta(x)]=\mathbb{E}[\beta(x+t)] \\
\mathbb{E}_{\beta \sim \mathcal{P}(\beta)}[\beta(x) \beta(y)]=\mathbb{E}[\beta(x+t) \beta(y+t)] \tag{3.3}
\end{array}
$$

i.e. first and second moments should be independent on the absolute location. The WienerKhintchine Theorem now states that the covariance has a spectral decomposition, i.e. it

[^1]is diagonal in Fourier space, under this condition (see e.g. Cha16]). Denoting the Fourier transform by $\mathcal{F}$, i.e. in the one dimensional case, $\mathcal{F}[f](q)=(2 \pi)^{-\frac{1}{2}} \int \mathrm{~d} x e^{-i q x} f(x)$. Therefore, the covariance can be completely specified by a one dimensional function:
\[

$$
\begin{equation*}
\left(\mathcal{F} B \mathcal{F}^{-1}\right)(k, q)=\delta(k-q) P_{\beta}(k) \tag{3.4}
\end{equation*}
$$

\]

Here, $P_{\beta}(k)$ is called the power spectrum in the formalism of IFT EFK09.

## A Poisson Lognormal Approach for Handling Discretization

Building on these considerations we now regard the problem of discretization. Measurement data itself is usually not purely continuous but can only be given in a somewhat discretized way (e.g. by the measurement device itself or by precision restrictions imposed from storing the data). Another problem is that many numerical approaches to inference tasks, such as Gaussian Process regression, use finite bases as approximations in order to efficiently obtain results Moo +16 Ste+10. Here, we aim to directly confront these problems by imposing a formalism where the discretization is inherent.

So instead of taking a direct approach with the above formulation, we use a Poissonian approach and consider an equidistant grid $\left\{z_{1}, \ldots, z_{n_{\text {bins }}}\right\}$ in the $[0,1]$ interval. This is equivalent to defining bins, where the $z_{j}$ are the midpoints of the bins. We now take the measurement counts, $k_{i}$ which gives the number of $x$-measurements within the $i$-th bin. For these measurement counts we now take a Poisson lognormal distribution as an Ansatz, that is, we assume that the measurement counts for the bins are Poisson distributed, where the means follow a lognormal distribution. We can model this discretization by applying a response operator $R: \mathbb{R}^{[0,1]} \rightarrow \mathbb{R}^{n_{\text {bins }}}$ to the lognormal field. This is done in the most direct way via employing a delta distribution

$$
\begin{equation*}
R_{j x} \equiv \delta\left(x-z_{j}\right) \tag{3.5}
\end{equation*}
$$

In order to allow for a more compact notation we will use an index notation from now on, e.g. $f_{x}=f(x)$ for some function $f$ or $O_{x y}=O(x, y)$ for some operator $O$. Whenever the indices are suppressed, an integration (in the continuous case) or dot product (in the discrete case) is understood, e.g. $(O f)_{x} \equiv O_{x y} f_{y}=\int \mathrm{d} y O_{x y} f_{y}=\int \mathrm{d} y O(x, y) f(y)$ In the following we will use bold characters for finite dimensional vectors, e.g. $\boldsymbol{\lambda} \equiv\left(\lambda_{1}, \ldots, \lambda_{n_{\text {bins }}}\right)^{T}$. By inserting such a finite dimensional vector in the argument of a function, e.g. $\beta(\boldsymbol{x})$ we refer to a vector consisting of the function evaluated at each entry of $\boldsymbol{x}$, that is $\left(\beta(\boldsymbol{z}) \equiv\left(\beta\left(z_{1}\right), \ldots, \beta\left(z_{n_{\text {bins }}}\right)\right)\right.$. Later on we will use the notation $\hat{\cdot}$ which raises some vector to a diagonal matrix ( $\hat{\boldsymbol{x}}_{i j} \equiv \delta_{i j} x_{i}$ ). We will use this notation analogously for fields, e.g. $\left(\hat{\beta}_{u v} \equiv \delta(u-v) \beta(u)\right)$. Now we can state
the probability distribution for the measurement counts $k_{j}$ :

$$
\begin{align*}
& \mathcal{P}\left(k_{j} \mid \lambda_{j}\right)=\frac{\lambda_{j}^{k_{j}} e^{-\lambda_{j}}}{k_{j}!}  \tag{3.7}\\
& \lambda_{j}=\mathbb{E}_{(k \mid \beta)}\left[k_{j}\right]=\rho e^{\beta_{z_{j}}}=\int d x R_{j x} e^{\beta_{x}}=\rho\left(R e^{\beta}\right)_{j}  \tag{3.8}\\
& \boldsymbol{\lambda}=\rho \boldsymbol{R} e^{\beta}=\rho e^{\beta(\boldsymbol{z})}  \tag{3.9}\\
& \mathcal{P}(\boldsymbol{k} \mid \boldsymbol{\lambda})=\prod_{j} \frac{\lambda_{j}^{k_{j}} e^{-\lambda_{j}}}{k_{j}!}=\prod_{j} \frac{\left(R_{j} e^{\beta}\right)^{k_{j}} e^{-R_{j} e^{\beta}}}{k_{j}!}=\frac{\left(\prod_{j}\left(R_{j} e^{\beta}\right)^{k_{j}}\right) e^{-\mathbf{1}^{\dagger} \boldsymbol{R} e^{\beta}}}{\prod_{j} k_{j}!}  \tag{3.10}\\
& \mathcal{P}(\boldsymbol{x} \mid \boldsymbol{k})=\frac{1}{N!} \tag{3.11}
\end{align*}
$$

Eq. 3.11 follows from the consideration that given the counts $\left(k_{1}, \ldots, k_{n_{\text {bins }}}\right)$ for the bins, only the positions of the observations $\left(x_{1}, \ldots, x_{N}\right)$ is fixed, but the ordering is not. The $N$ observations can be ordered in $N$ ! ways.

Now considering the whole $r$-vector of bin counts $\boldsymbol{k}$ at once, we get

$$
\begin{equation*}
\mathcal{P}(\boldsymbol{k} \mid \beta)=\frac{e^{\sum_{j} k_{j} \beta\left(z_{j}\right)} e^{-\rho^{\dagger} e^{\beta(z)}}}{\prod_{j} k_{j}!}=\frac{e^{\boldsymbol{k}^{\dagger} \beta(\boldsymbol{z})-\boldsymbol{\rho}^{\dagger} e^{\beta(z)}}}{\prod_{j} k_{j}!} \tag{3.12}
\end{equation*}
$$

Marginalizing $\beta$ we get

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{x} \mid P_{\beta}, X \rightarrow Y\right) & =\frac{1}{N!} \int \mathcal{D}[\beta] \mathcal{P}(x \mid \beta, X \rightarrow Y) \mathcal{P}\left(\beta \mid P_{\beta}\right) \\
& =\frac{1}{N!}|2 \pi B|^{-\frac{1}{2}} \int \mathcal{D}[\beta] \frac{e^{\boldsymbol{k}^{\dagger} \beta(\boldsymbol{z})-\boldsymbol{\rho}^{\dagger} e^{\beta(z)}}}{\prod_{j} k_{j}!} e^{-\frac{1}{2} \beta^{\dagger} B^{-1} \beta}= \\
& =\frac{|2 \pi B|^{-\frac{1}{2}}}{N!\prod_{j} k_{j}!} \int \mathcal{D}[\beta] e^{-\gamma[\beta]} \tag{3.14}
\end{align*}
$$

where

$$
\begin{equation*}
\gamma[\beta] \equiv-\boldsymbol{k}^{\dagger} \beta(\boldsymbol{z})+\boldsymbol{\rho}^{\dagger} e^{\beta(\boldsymbol{z})}+\frac{1}{2} \beta^{\dagger} B^{-1} \beta . \tag{3.15}
\end{equation*}
$$

We approach this integration by a saddle point approximation. In the following we will denote the functional derivative by $\partial$, i.e. $\partial_{f_{z}} \equiv \frac{\delta}{\delta f(z)}$.

Taking the first and second order functional derivative of $\gamma$ w.r.t. $\beta$ we get (

$$
\begin{align*}
& \partial_{\beta} \gamma[\beta]=-\boldsymbol{k}^{\dagger}+\rho\left(e^{\beta(\boldsymbol{z})}\right)^{\dagger}+\beta^{\dagger} B^{-1}  \tag{3.16}\\
& \partial_{\beta} \partial_{\beta} \gamma[\beta]=\widehat{\rho e^{\beta(\boldsymbol{z})}}+B^{-1} \tag{3.17}
\end{align*}
$$

The above derivatives are still defined in the space of functions $\mathbb{R}^{[0,1]}$, that is

$$
\begin{aligned}
& \boldsymbol{k}_{u}^{\dagger} \equiv \sum_{j=1}^{n_{\text {bins }}} k_{j}\left(\tilde{R}_{j}\right)_{u} \\
& \left(\widehat{\rho e^{\beta(z)}}\right)_{u v}=\rho \sum_{j=1}^{n_{\text {bins }}}\left(\tilde{R}_{j}\right)_{u}\left(\tilde{R}_{j}\right)_{v} e^{\beta(u)}
\end{aligned}
$$

i.e. a diagonal operator with $e^{\beta(z)}$ as diagonal entries.

Let $\beta_{0}$ denote the function that minimizes the functional $\gamma$, i.e.

$$
\begin{equation*}
\left.\frac{\delta \gamma[\beta]}{\delta \beta}\right|_{\beta=\beta_{0}}=0 \tag{3.18}
\end{equation*}
$$

We expand the functional $\gamma$ up to second order around $\beta_{0}$ :

$$
\begin{align*}
\int \mathcal{D}[\beta] e^{-\gamma[\beta]} & =\int \mathcal{D}[\beta] e^{-\gamma\left[\beta_{0}\right]-\left(\left.\frac{\delta \gamma[\beta]}{\delta \beta} \right\rvert\, \beta=\beta_{0}\right)^{\dagger} \beta-\frac{1}{2} \beta^{\dagger}\left(\left.\frac{\delta^{2} \gamma[\beta]}{\delta \beta^{\dagger} \beta} \right\rvert\, \beta=\beta_{0}\right) \beta+\mathcal{O}\left(\beta^{3}\right)} \\
& \approx e^{-\gamma\left[\beta_{0}\right]}\left|2 \pi\left(\left.\frac{\delta^{2} \gamma[\beta]}{\delta \beta^{2}}\right|_{\beta=\beta_{0}}\right)^{-1}\right|^{\frac{1}{2}} \\
& =e^{+\boldsymbol{k}^{\dagger} \boldsymbol{\beta}_{0}-\boldsymbol{\rho}^{\dagger} e^{\beta_{0}-\frac{1}{2} \beta_{0}^{\dagger} B^{-1} \beta_{0}}\left|\frac{1}{2 \pi}\left(\widehat{\left(\rho e^{\boldsymbol{\beta}_{0}}\right.}+B^{-1}\right)\right|^{-\frac{1}{2}}} \tag{3.19}
\end{align*}
$$

where we dropped higher order terms of $\beta$, used that the gradient at $\beta=\beta_{0}$ vanishes and evaluated the remaining Gaussian integral.

Plugging the result (3.19) into (3.14) and using

$$
\begin{equation*}
|2 \pi B|^{-\frac{1}{2}}\left|\frac{1}{2 \pi}\left(\widehat{\rho e^{\boldsymbol{\beta}_{\mathbf{0}}}}+B^{-1}\right)\right|^{-\frac{1}{2}}=\left|B\left(\widehat{\rho e^{\boldsymbol{\beta}_{\mathbf{0}}}}+B^{-1}\right)\right|^{-\frac{1}{2}}=\left|\rho B \widehat{e^{\boldsymbol{\beta}_{\mathbf{0}}}}+\mathbb{1}\right|^{-\frac{1}{2}} \tag{3.20}
\end{equation*}
$$

we get:

$$
\begin{align*}
& \mathcal{P}\left(\boldsymbol{x} \mid P_{\beta}, X \rightarrow Y\right) \approx \frac{1}{N!} \frac{e^{+\boldsymbol{k}^{\dagger} \boldsymbol{\beta}_{\mathbf{0}}-\boldsymbol{\rho}^{\dagger} e^{\boldsymbol{\beta}_{\mathbf{0}}}-\frac{1}{2} \beta_{0}^{\dagger} B^{-1} \beta_{0}}}{\left|\rho B \widehat{\boldsymbol{\beta}_{\mathbf{0}}}+\mathbb{1}\right|^{\frac{1}{2}} \prod_{j} k_{j}!}  \tag{3.21}\\
& \mathcal{H}\left(\boldsymbol{x} \mid P_{\beta}, X \rightarrow Y\right) \approx \mathcal{H}_{0}+\frac{1}{2} \log \left|\rho B \widehat{e^{\boldsymbol{\beta}_{\mathbf{0}}}}+\hat{\mathbb{1}}\right|+\log \left(\prod_{j} k_{j}!\right)-\boldsymbol{k}^{\dagger} \boldsymbol{\beta}_{\mathbf{0}}+\boldsymbol{\rho}^{\dagger} e^{\boldsymbol{\beta}_{\mathbf{0}}}+\frac{1}{2} \beta_{0}^{\dagger} B^{-1} \beta_{0} \tag{3.22}
\end{align*}
$$

where $\mathcal{H}(\cdot) \equiv-\log (\mathcal{P}(\cdot))$ is called the information Hamiltonian and $\mathcal{H}_{0}$ collects all terms which do not depend on the data $\boldsymbol{d}$.

### 3.2.2. Functional Relation of Cause and Effect

Similarly to $\beta$, we suppose a Gaussian distribution for the function $f$, relating $Y$ to $X$ :

$$
\begin{equation*}
\mathbb{R}^{[0,1]} \ni f \sim \mathcal{N}(0 \mid f, F) \tag{3.23}
\end{equation*}
$$

Proposing a Fourier diagonal covariance $F$ once more, determined by a power spectrum $P_{f}$ :

$$
\begin{equation*}
\left(\mathcal{F} F \mathcal{F}^{-1}\right)(k, q)=\delta(k-q) P_{f}(k) \tag{3.24}
\end{equation*}
$$

we assume additive Gaussian noise, using the notation $f(\boldsymbol{x}) \equiv\left(f\left(x_{1}\right), \ldots, f\left(x_{N}\right)\right)^{T}$ and $\boldsymbol{\epsilon} \equiv$ $\left(\epsilon_{1}, \ldots, \epsilon_{N}\right)^{T}$, we have

$$
\begin{align*}
& \boldsymbol{y}=f(\boldsymbol{x})+\boldsymbol{\epsilon}  \tag{3.25}\\
& \boldsymbol{\epsilon} \sim \mathcal{N}(\epsilon \mid 0, \mathcal{E})  \tag{3.26}\\
& \mathcal{E} \equiv \operatorname{diag}\left(\varsigma^{2}, \varsigma^{2}, \ldots\right)=\varsigma^{2} \mathbb{1} \in \mathbb{R}^{N \times N} \tag{3.27}
\end{align*}
$$

that is each independent noise sample is drawn from a zero-mean Gaussian distribution with given variance $\varsigma^{2}$.

Knowing the noise $\boldsymbol{e}$, the cause $\boldsymbol{x}$ and the causal mechanism $f$ completely determines $\boldsymbol{y}$ via 3.25. Therefore, $\mathcal{P}(\boldsymbol{y} \mid \boldsymbol{x}, f, \boldsymbol{\epsilon}, X \rightarrow Y)=\delta(\boldsymbol{y}-f(\boldsymbol{x})-\boldsymbol{\epsilon})$, where $\delta(\cdot)$ denotes the Dirac delta distribution. We can now state the conditional distribution for the effect variable measurements $\boldsymbol{y}$, given the cause variable measurements $\boldsymbol{x}$. Marginalizing out the dependence on the relating function $f$ and the noise $\boldsymbol{\epsilon}$ we get:

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, P_{f}, \varsigma, X \rightarrow Y\right) & =\int \mathcal{D}[f] \mathrm{d}^{N} \boldsymbol{\epsilon} \mathcal{P}(\boldsymbol{y} \mid \boldsymbol{x}, f, \boldsymbol{\epsilon}, X \rightarrow Y) \mathcal{P}(\boldsymbol{\epsilon} \mid \varsigma) \mathcal{P}\left(f \mid P_{f}\right) \\
& =\int \mathcal{D}[f] \mathrm{d}^{N} \boldsymbol{\epsilon} \delta(\boldsymbol{y}-f(\boldsymbol{x})-\boldsymbol{\epsilon}) \mathcal{N}(\boldsymbol{\epsilon} \mid 0, \mathcal{E}) \mathcal{N}(f \mid 0, F) \tag{3.28}
\end{align*}
$$

We will now use the Fourier representation of the delta distribution, specifically $\delta(x)=$ $\int \frac{\mathrm{d} q}{2 \pi} e^{i q x}$

$$
\begin{equation*}
\delta(\boldsymbol{y}-f(\boldsymbol{x})-\boldsymbol{\epsilon})=\int \frac{\mathrm{d}^{N} \boldsymbol{q}}{(2 \pi)^{N}} e^{i \boldsymbol{q}^{\dagger}(\boldsymbol{y}-\boldsymbol{\epsilon}-f(\boldsymbol{x}))}=\int \frac{\mathrm{d}^{N} \boldsymbol{q}}{(2 \pi)^{N}} e^{i \boldsymbol{q}^{\dagger}(\boldsymbol{y}-\boldsymbol{\epsilon}-f(\boldsymbol{x}))} \tag{3.29}
\end{equation*}
$$

Once more we employ a vector of response operators, mapping $\mathbb{R}^{\mathbb{R}}$ to $\mathbb{R}^{N}$,

$$
\begin{equation*}
\boldsymbol{R}_{x} \equiv\left(R_{1 x}, \ldots, R_{N x}\right)^{T}=\left(\delta\left(x-x_{1}\right), \ldots, \delta\left(x-x_{N}\right)\right)^{T} \tag{3.30}
\end{equation*}
$$

This allows to represent the evaluation $f(\boldsymbol{x})=\boldsymbol{R} \dagger f$, i.e. as a linear dot-product. Using the well known result for Gaussian integrals with linear terms (see e.g. GBR13),

$$
\begin{equation*}
\int \mathcal{D}[u] e^{-\frac{1}{2} u^{\dagger} A u+b^{\dagger} u}=\left|\frac{A}{2 \pi}\right|^{-\frac{1}{2}} e^{\frac{1}{2} b^{\dagger} A b} \tag{3.31}
\end{equation*}
$$

we are able to analytically do the path integral over $f$,

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, P_{f}, \varsigma, X \rightarrow Y\right) & =|2 \pi F|^{-\frac{1}{2}} \int \mathcal{D}[f] \mathrm{d}^{N} \boldsymbol{\epsilon} \frac{\mathrm{~d}^{N} \boldsymbol{q}}{(2 \pi)^{N}} e^{i \boldsymbol{q}^{\dagger}\left(\boldsymbol{y}-\boldsymbol{\epsilon}-\boldsymbol{R}^{\dagger} f\right)-\frac{1}{2} f^{\dagger} F^{-1} f} \mathcal{N}(\boldsymbol{\epsilon} \mid 0, \mathcal{E}) \\
& =\int \mathrm{d}^{N} \boldsymbol{\epsilon} \frac{\mathrm{~d}^{N} \boldsymbol{q}}{(2 \pi)^{N}} e^{i \boldsymbol{q}^{\dagger}(\boldsymbol{y}-\boldsymbol{\epsilon})+(-i)^{2} \frac{1}{2} \boldsymbol{q}^{\dagger} \boldsymbol{R}^{\dagger} F \boldsymbol{R} \boldsymbol{q}} \mathcal{N}(\boldsymbol{\epsilon} \mid 0, \mathcal{E}) \tag{3.32}
\end{align*}
$$

Now we do the integration over the noise variable, $\boldsymbol{\epsilon}$, by using the equivalent of Eq. 3.31 for the vector-valued case:

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, P_{f}, \varsigma, X \rightarrow Y\right) & =|2 \pi \mathcal{E}|^{-\frac{1}{2}} \int \mathrm{~d}^{N} \boldsymbol{\epsilon} \frac{\mathrm{~d}^{N} \boldsymbol{q}}{(2 \pi)^{N}} e^{i \boldsymbol{q}^{\dagger}(\boldsymbol{y}-\boldsymbol{\epsilon})-\frac{1}{2} \boldsymbol{q}^{\dagger} \boldsymbol{R}^{\dagger} F \boldsymbol{R} \boldsymbol{q}-\frac{1}{2} \epsilon^{\dagger} \mathcal{E}^{-1} \boldsymbol{\epsilon}} \\
& =\int \frac{\mathrm{d}^{N} \boldsymbol{q}}{(2 \pi)^{N}} e^{i \boldsymbol{q}^{\dagger} \boldsymbol{y}-\frac{1}{2} \boldsymbol{q}\left(\boldsymbol{R}^{\dagger} F \boldsymbol{R}+\mathcal{E}\right) \boldsymbol{q}} \tag{3.33}
\end{align*}
$$

In the following we will write $\mathbb{R}^{N \times N} \ni \tilde{F}=\boldsymbol{R}^{\dagger} F \boldsymbol{R}$, with entries $\tilde{F}_{i j}=F\left(x_{i}, x_{j}\right) .^{2}$ The integration over the Fourier modes $\boldsymbol{q}$, again via the multivariate equivalent of 3.31, will give the preliminary result:

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, P_{f}, \varsigma, X \rightarrow Y\right) & =\int \frac{\mathrm{d}^{N} \boldsymbol{q}}{(2 \pi)^{N}} e^{i \boldsymbol{\boldsymbol { q } ^ { \dagger } \boldsymbol { y } - \frac { 1 } { 2 } \boldsymbol { q } ( \tilde { F } + \mathcal { E } ) \boldsymbol { q }}} \\
& =(2 \pi)^{-\frac{N}{2}}|\tilde{F}+\mathcal{E}|^{-\frac{1}{2}} e^{-\frac{1}{2} \boldsymbol{y}^{\dagger}(\tilde{F}+\mathcal{E})^{-1} \boldsymbol{y}} \tag{3.34}
\end{align*}
$$

### 3.2.3. Computing the Bayes factor

Now we are able to calculate the full likelihood of the data $\boldsymbol{d}=(\boldsymbol{x}, \boldsymbol{y})$ given our assumptions $P_{\beta}, P_{f}, \varsigma$ for the direction $X \rightarrow Y$ and vice versa $Y \rightarrow X$. As we are only interested in the ratio of the probabilities and not in the absolute probabilities itself, it suffices to calculate the Bayes factor:

$$
\begin{align*}
O_{X \rightarrow Y} & =\frac{\mathcal{P}\left(\boldsymbol{d} \mid P_{\beta}, P_{f}, \varsigma, X \rightarrow Y\right)}{\mathcal{P}\left(\boldsymbol{d} \mid P_{\beta}, P_{f}, \varsigma, Y \rightarrow X\right)} \\
& =\exp \left[\mathcal{H}\left(\boldsymbol{d} \mid P_{\beta}, P_{f}, \varsigma, Y \rightarrow X\right)-\mathcal{H}\left(\boldsymbol{d} \mid P_{\beta}, P_{f}, \varsigma, X \rightarrow Y\right)\right] \tag{3.35}
\end{align*}
$$

Above we used again the information Hamiltonian $\mathcal{H}(\cdot) \equiv-\log \mathcal{P}(\cdot)$
Making use of (3.21) and (3.34) we get, using the calculus for conditional distributions on the Hamiltonians, $\mathcal{H}(A, B)=\overline{\mathcal{H}}(A \mid B)+\mathcal{H}(B)$

$$
\left.\begin{align*}
\mathcal{H}\left(\boldsymbol{d} \mid P_{\beta}, P_{f}, \varsigma, X \rightarrow Y\right) & =\mathcal{H}\left(\boldsymbol{x} \mid P_{\beta}, X \rightarrow Y\right)+\mathcal{H}\left(\boldsymbol{y} \mid \boldsymbol{x}, P_{f}, \varsigma, X \rightarrow Y\right) \\
& \left.=\mathcal{H}_{0}+\log \left(\prod_{j} k_{j}!\right)+\frac{1}{2} \log \right\rvert\, \rho B e^{\boldsymbol{\beta}_{\mathbf{0}}}
\end{align*}+\mathbb{1} \right\rvert\,-\boldsymbol{k}^{\dagger} \boldsymbol{\beta}_{\mathbf{0}}+\quad .
$$

[^2]Where we suppressed the dependence of $\tilde{F}, \beta_{0}$ on $\boldsymbol{x}$ (for the latter, the dependence is not explicit, but rather implicit as $\beta_{0}$ is determined by the minimum of the $\boldsymbol{x}$-dependent functional $\gamma)$.

We omit stating $\mathcal{H}\left(\boldsymbol{d} \mid P_{\beta}, P_{f}, \varsigma, Y \rightarrow X\right)$ explicitly as the expression is just given by taking (3.36) and switching $\boldsymbol{x}$ and $\boldsymbol{y}$ or $X$ and $Y$, respectively.

### 3.3. Inference of the Noise Variance

### 3.3.1. Unknown Variance of the Noise

So far we assumed the variance of the noise to be known and identical, $\epsilon_{i} \sim \mathcal{N}\left(\epsilon_{i} \mid 0, \varsigma^{2}\right)$. We want to relax this condition, allowing the noise variance to be a priori unknown and position dependent by substituting $\varsigma^{2}=\varsigma(x)^{2} \rightarrow e^{\eta(x)}$, where $\eta \in \mathbb{R}^{[0,1]}$

$$
\begin{align*}
& \boldsymbol{\epsilon} \sim \mathcal{N}(\boldsymbol{\epsilon} \mid 0, \mathcal{E}(\boldsymbol{x}))  \tag{3.37}\\
& \mathcal{E}(\boldsymbol{x})=\operatorname{diag}\left(e^{\eta\left(x_{1}\right)}, e^{\eta\left(x_{2}\right)}, \ldots, e^{\eta\left(x_{N}\right)}\right)=\widehat{e^{\eta(\boldsymbol{x})}}  \tag{3.38}\\
& \Rightarrow|\mathcal{E}|=\prod_{i}^{N} e^{\eta\left(x_{i}\right)}=e^{\mathbf{1}^{\dagger} \eta(\boldsymbol{x})} \tag{3.39}
\end{align*}
$$

We further argue that there exists a spatial correlation in the noise level, meaning that points $x, x^{\prime}$ which are close are affected by a similar noise variance $e^{\eta(x)}, e^{\eta\left(x^{\prime}\right)}$. We incorporate this consideration by constraining the derivative $\nabla e^{\eta(x)}$ with a prior:

$$
\begin{equation*}
\mathcal{P}\left(\eta \mid \sigma_{\eta}\right)=Z\left[\sigma_{\eta}\right]^{-1} e^{-\frac{1}{2 \sigma_{\eta}} \int \mathrm{d} q|\nabla \eta(q)|^{2}} \propto e^{-\frac{1}{2 \sigma_{\eta}} \eta^{\dagger} \nabla^{\dagger} \nabla \eta} \tag{3.40}
\end{equation*}
$$

where we dropped the normalization factor $Z\left[\sigma_{\eta}\right]=\int \mathcal{D}[\eta] e^{-\frac{1}{2 \sigma_{\eta}} \eta^{\dagger} \nabla^{\dagger} \nabla \eta}=\left|\frac{\nabla^{\dagger} \nabla}{2 \pi \sigma_{\eta}}\right|^{-\frac{1}{2}}$ as it depends only on the hyper parameter $\sigma_{\eta}$ which we can set to a canonical value.

We remark that it would be possible to introduce a similar inference structure for the function $\eta$ as for the functions $\beta$ and $f$. This would mean a functional normal distribution for $\eta$ with a Fourier diagonal covariance matrix given by a power spectrum $P_{\eta}$ which itself can be assumed to be set as a fixed hyperparameter (or to be just constrained by some prior distribution and made part of the inference). However at this point we choose to refrain from this step in order to avoid the resulting model from being unnecessary complex.


Figure 3.2.: Overview over the Bayesian hierarchical model, for the case $X \rightarrow Y$. The power spectra $P_{\beta}, P_{f}$ are given as hyperparameters, but the noise variance is part of the inference.

### 3.3.2. Marginalization of the Noise Variance

Marginalizing out the noise variance in (3.34) gives

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, P_{f}, \sigma_{\eta}, X \rightarrow Y\right) & =\int \mathcal{D}[f, \eta] \mathrm{d}^{N} \boldsymbol{\epsilon} \mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, \tau_{f}, \eta, X \rightarrow Y\right) \mathcal{P}(\boldsymbol{\epsilon} \mid \eta) \mathcal{P}\left(f \mid P_{f}\right) \mathcal{P}\left(\eta \mid \sigma_{\eta}\right) \propto \\
& \propto \int \mathcal{D}[\eta]\left|\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right|^{-\frac{1}{2}} \exp \left(-\frac{1}{2} \boldsymbol{y}^{\dagger}\left(\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right)^{-1} \boldsymbol{y}-\frac{1}{2 \sigma_{\eta}} \eta \nabla^{\dagger} \nabla \eta\right) \tag{3.41}
\end{align*}
$$

again omitting the normalization factors depending only on $\sigma_{\eta}$ and the $\frac{1}{\sqrt{2 \pi}^{N}}$ factor.
In order to tackle the integration we again employ the Laplace approximation, introducing the energy functional $\gamma_{\eta}: \mathbb{R}^{[0,1]} \rightarrow \mathbb{R}$ :

$$
\begin{equation*}
\gamma_{\eta}[\eta] \equiv \frac{1}{2} \log \left|\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right|+\frac{1}{2} \boldsymbol{y}^{\dagger}\left(\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right)^{-1} \boldsymbol{y}+\frac{1}{2 \sigma_{\eta}} \eta \nabla^{\dagger} \nabla \eta \tag{3.42}
\end{equation*}
$$

We now need to compute the first and second order derivatives of $\gamma_{\eta}(\eta)$ w.r.t. $\eta$. Especially the derivatives of the logarithmic determinant $\log \left(\left|\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right|\right)$ are non-obvious, so we will
explicitly derive these. We begin by remarking that for some non-singular matrix or operator $A$, we have the identity [SN10]

$$
\begin{equation*}
\log |A|=\operatorname{tr} \log (A) \tag{3.43}
\end{equation*}
$$

Applying the chain rule here gives the formula:

$$
\begin{equation*}
\frac{\partial|A[t]|}{\partial t}=|A[t]| \operatorname{tr}\left[A^{-1}\left(\frac{\partial A[t]}{\partial t}\right)\right] \tag{3.44}
\end{equation*}
$$

We further make the abbreviation:

$$
\begin{equation*}
G[\eta] \equiv\left(\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right)^{-1} \tag{3.45}
\end{equation*}
$$

Finally we will use the relation:

$$
\begin{align*}
0=\frac{\partial \mathbb{1}}{\partial t} & =\frac{\partial A[t] A[t]^{-1}}{\partial t}=\frac{\partial A[t]}{\partial t} A[t]^{-1}+A[t] \frac{\partial A[t]^{-1}}{\partial t} \\
\Rightarrow \frac{\partial A[t]^{-1}}{\partial t} & =-A[t]^{-1} \frac{\partial A[t]}{\partial t} A[t]^{-1} \tag{3.46}
\end{align*}
$$

We therefore get:

$$
\begin{align*}
\partial_{\eta_{u}} \log \left|\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right| & =\operatorname{tr}\left(G[\eta] \partial_{\eta_{u}}\left(\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right)\right)=\operatorname{tr}\left(G[\eta] \widehat{\left.e^{\widehat{\eta(\boldsymbol{x}) \delta_{\boldsymbol{x}}}}\right)}\right.  \tag{3.47}\\
\partial_{\eta_{v}} \partial_{\eta_{u}} \log \left|\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right| & =\operatorname{tr}\left(\left(\partial_{\eta_{v}} G\right) \widehat{e^{\eta(\boldsymbol{x}) \delta_{\boldsymbol{x}}}}+G[\eta]\left(\partial_{\eta_{v}} \widehat{\left.e^{\eta \boldsymbol{x}) \delta_{\boldsymbol{x} u}}\right)}\right)\right. \\
& =\operatorname{tr}\left(-G[\eta] \widehat{\left.e^{\eta(\boldsymbol{x}) \delta_{\boldsymbol{x} v}} G[\eta] \widehat{e^{\eta(\boldsymbol{x}) \delta_{\boldsymbol{x}}}}+\delta_{u v} G[\eta] \widehat{e^{\eta(\boldsymbol{x}) \delta_{\boldsymbol{x}} u}}\right)}\right. \tag{3.48}
\end{align*}
$$

The derivatives of the other terms are rather straight-forward:

$$
\begin{align*}
\partial_{\eta_{u}} \boldsymbol{y}^{\dagger}\left(\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right)^{-1} \boldsymbol{y} & =-\boldsymbol{y}^{\dagger} G[\eta] \widehat{e^{\eta(\boldsymbol{x}) \delta_{\boldsymbol{x}}} G[\eta] \boldsymbol{y}}  \tag{3.49}\\
\partial_{\eta_{v}} \partial_{\eta_{u}} \boldsymbol{y}^{\dagger}\left(\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right)^{-1} \boldsymbol{y} & =\boldsymbol{y}^{\dagger}\left(2 G[\eta] \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G[\eta] \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} v}} G[\eta]-G[\eta] \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} \delta_{u v} G[\eta]\right) \boldsymbol{y}  \tag{3.50}\\
\partial_{\eta_{u}} \frac{1}{2 \sigma_{\eta}} \eta \nabla^{\dagger} \nabla \eta & =\frac{1}{\sigma_{\eta}}\left(\eta \nabla^{\dagger} \nabla\right)_{u} \\
\partial_{\eta_{v}} \partial_{\eta_{u}} \frac{1}{2 \sigma_{\eta}} \eta \nabla^{\dagger} \nabla \eta & =\frac{1}{\sigma_{\eta}}\left(\nabla^{\dagger} \nabla\right)_{u v} \tag{3.51}
\end{align*}
$$

We thus have the gradient:

$$
\begin{equation*}
\partial_{\eta_{u}} \gamma_{\eta}[\eta]=\frac{1}{2} \operatorname{tr}\left(G[\eta] \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}}\right)-\frac{1}{2} \boldsymbol{y}^{\dagger} G[\eta] \widehat{e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}} u} G[\eta] \boldsymbol{y}+\frac{1}{\sigma_{\eta}}\left(\eta \nabla^{\dagger} \nabla\right)_{u} .{ }_{u} .} \tag{3.52}
\end{equation*}
$$

and the curvature

$$
\begin{align*}
\partial_{\eta_{v}} \partial_{\eta_{u}} \gamma_{\eta}[\eta]= & \operatorname{tr}\left(-G[\eta] \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} v}} G[\eta] \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}}}+\delta_{u v} G[\eta] \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}}}\right) \\
& +\frac{1}{2} \boldsymbol{y}^{\dagger}\left(2 G[\eta] \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G[\eta] \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} v}} G[\eta]-G[\eta] e^{\hat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} \delta_{u v} G[\eta]\right) \boldsymbol{y} \\
& +\frac{1}{\sigma_{\eta}}\left(\nabla^{\dagger} \nabla\right)_{u v} \equiv \\
\equiv & \Gamma_{\eta}[\eta] \tag{3.53}
\end{align*}
$$

We write the the density in 3.41 in an exponential form, using the energy defined in 3.42 .

$$
\begin{equation*}
\mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, P_{f}, \sigma_{\eta}, X \rightarrow Y\right) \propto \int \mathcal{D}[\eta] e^{-\gamma_{\eta}[\eta]} \tag{3.54}
\end{equation*}
$$

And expand up to second order in $\eta$ around $\eta_{0}=\underset{\eta \in \mathbb{R}^{[0,1]}}{\operatorname{argmin}} \gamma_{\eta}[\eta]$ :

$$
\begin{align*}
& \mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, P_{f}, \sigma_{\eta}, X \rightarrow Y\right) \propto \int \mathcal{D}[\eta] e^{-\gamma_{n}\left[\eta_{0}\right]-\left(\left.\frac{\delta \gamma_{[ }[\eta]}{\delta \eta} \right\rvert\, \eta=\eta_{0}\right)^{\dagger} \eta-\frac{1}{2} \eta^{\dagger}\left(\left.\frac{\delta^{2} \gamma_{\eta}[\eta]}{\delta \eta^{\dagger} \eta}\right|_{\eta=\eta_{0}}\right) \eta+\mathcal{O}\left(\eta^{3}\right)} \\
& \approx e^{-\gamma_{\eta}\left[\eta_{0}\right]}\left|2 \pi\left(\left.\frac{\delta^{2} \gamma_{\eta}[\eta]}{\delta \eta^{2}}\right|_{\eta=\eta_{0}}\right)^{-1}\right|^{\frac{1}{2}} \\
&=e^{\left.-\frac{1}{2} \log \right\rvert\, \tilde{F}+e^{\eta_{0}(x)}} \left\lvert\,-\frac{1}{2} \boldsymbol{y}^{\dagger}\left(\tilde{F}+e^{\eta_{0}(x)}\right)^{-1} \boldsymbol{y}-\frac{1}{2 \sigma_{\eta}} \eta_{0} \nabla^{\dagger} \nabla \eta_{0}\right.  \tag{3.55}\\
&\left|\frac{1}{2 \pi} \Gamma_{\eta}\left[\eta_{0}\right]\right|^{-\frac{1}{2}}
\end{align*}
$$

Using the result above, we can state the information Hamiltonian for the causal direction $X \rightarrow Y$ :

$$
\begin{align*}
\mathcal{H}\left(\boldsymbol{d} \mid P_{\beta}, P_{f}, \sigma_{\eta}, X \rightarrow Y\right) & =\mathcal{H}\left(\boldsymbol{x} \mid P_{\beta}, X \rightarrow Y\right)+\mathcal{H}\left(\boldsymbol{y} \mid \boldsymbol{x}, P_{f}, \sigma_{\eta}, X \rightarrow Y\right) \\
& =\mathcal{H}_{0}+\log \left(\prod_{j} k_{j}!\right)+\frac{1}{2} \log \left|\rho B \widehat{e^{\boldsymbol{\beta}_{0}}}+\mathbb{1}\right|-\boldsymbol{k}^{\dagger} \boldsymbol{\beta}_{\mathbf{0}}+\boldsymbol{\rho}^{\dagger} e^{\boldsymbol{\beta}_{\mathbf{0}}} \\
& +\frac{1}{2} \log \left|\frac{1}{2 \pi} \Gamma_{\eta}\left[\eta_{0}\right]\right|+\frac{1}{2} \beta_{0}^{\dagger} B^{-1} \beta_{0}+\frac{1}{2} \log \left|\tilde{F}+\widehat{e^{\eta_{0}(\boldsymbol{x} \boldsymbol{x}}}\right| \\
& +\frac{1}{2} \boldsymbol{y}^{\dagger}\left(\tilde{F}+\widehat{e^{\eta_{0}(\boldsymbol{x})}}\right)^{-1} \boldsymbol{y}+\frac{1}{2 \sigma_{\eta}} \eta_{0} \nabla^{\dagger} \nabla \eta_{0} \tag{3.56}
\end{align*}
$$

### 3.4. Inference of Power Spectra

### 3.4.1. Unknown Power Spectra

So far we assumed the power spectra $P_{\beta}, P_{f}$ to be given. We can approach a deeper model selection by expanding our inference on the power spectra itself and just assuming a somewhat
smooth structure for the power spectra as a prior distribution. Assuming smoothness on a logarithmic scale gives:

$$
\begin{align*}
& P_{\beta / f}(q)=e^{\tau_{\beta / f}(q)}  \tag{3.57}\\
& \mathcal{H}\left(\tau_{\beta / f} \mid \sigma_{\beta / f}\right)=-\log \left(Z\left[\sigma_{\beta / f}\right]\right) \frac{1}{2} \int \mathrm{~d} \log (q) \sigma_{\beta / f}^{-2}\left(\frac{\partial^{2} \tau_{\beta / f}(q)}{\partial(\log q)^{2}}\right)^{2}  \tag{3.58}\\
& \mathcal{P}\left(\tau_{\beta / f} \mid \sigma_{\beta / f}\right)=Z\left[\sigma_{\beta / f}\right]^{-1} e^{-\frac{1}{2 \sigma_{\beta / f}} \tau_{\beta / f}^{\dagger} \Delta^{\dagger} \Delta \tau_{\beta / f}} \tag{3.59}
\end{align*}
$$

here,

$$
\begin{equation*}
Z\left[\sigma_{\beta / f}\right] \equiv \int \mathcal{D}[\tau] e^{-\frac{1}{2 \sigma_{\beta / f}} \tau^{\dagger} \Delta^{\dagger} \Delta \tau}=\left|\frac{\Delta^{\dagger} \Delta}{2 \pi \sigma_{\beta / f}}\right|^{-\frac{1}{2}} \tag{3.60}
\end{equation*}
$$

serves as a normalization factor.

### 3.4.2. Marginalization of the Cause Distribution Power Spectrum

We can now perform a marginalization of the power spectrum $P_{\beta}=e^{\tau_{\beta}}$ in (3.14). We can use that $B=\left(\mathcal{F}^{\dagger} e^{\tau_{\beta}} \mathcal{F}\right)$ and thus $|B|=\left|\mathcal{F}^{\dagger}\right|\left|e^{\tau_{\beta}}\right||\mathcal{F}|=\left|\widehat{e}^{\tau_{\beta}}\right|$

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{x} \mid \sigma_{\beta}, X \rightarrow Y\right) & =\int \mathcal{D}\left[\tau_{\beta}\right] \mathcal{P}\left(\boldsymbol{x} \mid \tau_{\beta}, \sigma_{\beta}, X \rightarrow Y\right) \mathcal{P}\left(\tau_{\beta} \mid \sigma_{\beta}\right) \propto \\
& \propto \int \mathcal{D}\left[\beta, \tau_{\beta}\right] \frac{\left|2 \pi e^{\widehat{\tau_{\beta}}}\right|^{-\frac{1}{2}}}{\prod_{j} k_{j}!} e^{\boldsymbol{k}^{\dagger} \beta(\boldsymbol{z})-\boldsymbol{\rho}^{\dagger} e^{\beta(z)}-\frac{1}{2} \beta^{\dagger} \mathcal{F}^{\dagger} e^{-\tau_{\beta}} \mathcal{F} \beta-\frac{1}{2 \sigma_{\beta}} \tau_{\beta}^{\dagger} \Delta^{\dagger} \Delta \tau_{\beta}} \tag{3.61}
\end{align*}
$$

where we used $B^{-1}=\left(\mathcal{F}^{\dagger} \widehat{e}^{\tau_{\beta}} \mathcal{F}\right)^{-1}=\mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \mathcal{F}}$

### 3.4.3. Marginalization of the Power Spectrum for the Causal Mechanism

We perform a similar marginalization over $\tau_{f}$ for the likelihood of $\boldsymbol{y}$, given $\boldsymbol{x}$. We begin by noting:

$$
\begin{align*}
\tilde{F}_{i j} & =F\left(x_{i}, x_{j}\right)=\left(\mathcal{F}^{\dagger} e^{\tau_{f}} \mathcal{F}\right)\left(x_{i}, x_{j}\right) \\
& =\int \frac{1}{2 \pi} \mathrm{~d} q \mathrm{~d} q^{\prime} e^{i x_{i} q} e^{-i x_{j} q^{\prime}} e^{\tau_{f}(q)} \delta\left(q-q^{\prime}\right) \\
& =\int \frac{1}{2 \pi} \mathrm{~d} q e^{i q\left(x_{i}-x_{j}\right)} e^{\tau_{f}(q)} \\
& =\mathcal{F}\left[e^{\tau_{f}}\right]\left(x_{i}-x_{j}\right) \tag{3.62}
\end{align*}
$$

From (3.41) we now get

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, \sigma_{f}, \eta\right)_{X \rightarrow Y} & =\int \mathcal{D}\left[\tau_{f}\right] \mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, \tau_{f}, \eta\right)_{X \rightarrow Y} \mathcal{P}\left(\tau_{f} \mid \sigma_{f}\right) \propto \\
\propto & \left.\int \mathcal{D}\left[\tau_{f}\right] \widehat{e^{\eta(\boldsymbol{x})}}\right|^{\frac{1}{2}}\left|\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right|^{-\frac{1}{2}} e^{-\frac{1}{2} \boldsymbol{y}^{\dagger}\left(\tilde{F}+\widehat{\left.e^{\eta(\boldsymbol{x})}\right)^{-1}} \boldsymbol{y}-\frac{1}{2 \sigma_{f}} \tau_{f}^{\dagger} \Delta^{\dagger} \Delta \tau_{f}\right.}  \tag{3.63}\\
\mathcal{H}\left(\boldsymbol{y}, \tau_{f} \mid \boldsymbol{x}, \sigma_{f}, \eta\right)_{X \rightarrow Y}+\mathcal{H}_{1} & -\frac{1}{2} \log \left|\widehat{e^{\eta(\boldsymbol{x})}}\right|+\frac{1}{2} \log \left|\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right| \\
& +\frac{1}{2} \boldsymbol{y}^{\dagger}\left(\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right)^{-1} \boldsymbol{y}-\frac{1}{2 \sigma_{f}} \tau_{f}^{\dagger} \Delta^{\dagger} \Delta \tau_{f} \tag{3.64}
\end{align*}
$$

where we again collected data-independent factors as $\frac{1}{\sqrt{2 \pi}^{N}}$ and $Z\left[\sigma_{f}\right]^{-1}$ in $\mathcal{H}_{1}$.
This leads to the marginalized likelihood for $\boldsymbol{y}$ given $\boldsymbol{x}$ for the deep model,

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, P_{f}, \sigma_{\eta}, X \rightarrow Y\right) & =\int \mathcal{D}[f, \eta] \mathrm{d}^{N} \boldsymbol{\epsilon} \mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, \tau_{f}, \eta, X \rightarrow Y\right) \mathcal{P}(\boldsymbol{\epsilon} \mid \eta) \mathcal{P}\left(f \mid \sigma_{f}\right) \mathcal{P}\left(\eta \mid \sigma_{\eta}\right) \propto \\
& \propto \int \mathcal{D}[\eta]\left|\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right|^{-\frac{1}{2}} \exp \left(-\frac{1}{2} \boldsymbol{y}^{\dagger}\left(\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right)^{-1} \boldsymbol{y}+\frac{1}{2} \mathbf{1}^{\dagger} \eta(\boldsymbol{x})\right. \\
& \left.-\frac{1}{2 \sigma_{f}} \tau_{f}^{\dagger} \Delta^{\dagger} \Delta \tau_{f}-\frac{1}{2 \sigma_{\eta}} \eta \nabla^{\dagger} \nabla \eta\right) \tag{3.65}
\end{align*}
$$

### 3.4.4. Evidence for the Deep Hierarchical Model

We summarize the background assumptions (which includes the model, the hyper parameters and the causal direction) into a hypothesis $H_{1}$ and $H_{-1}$ respectively:

$$
\begin{align*}
H_{1} & \equiv\left(\sigma_{\beta}, \sigma_{f}, \sigma_{\eta}, X \rightarrow Y\right) \\
H_{-1} & \equiv\left(\sigma_{\beta}, \sigma_{f}, \sigma_{\eta}, Y \rightarrow X\right) \tag{3.66}
\end{align*}
$$

For the full evidence, given the Hypothesis $H_{1}$ :

$$
\begin{align*}
& \mathcal{P}\left(\boldsymbol{d} \mid H_{1}\right)= \mathcal{P}\left(\boldsymbol{x} \mid \sigma_{\beta}, X \rightarrow Y\right) \mathcal{P}\left(\boldsymbol{y} \mid \boldsymbol{x}, \sigma_{f}, \sigma \eta, X \rightarrow Y\right) \propto \\
& \propto\left.\int \mathcal{D}\left[\beta, \tau_{\beta}, \tau_{f}, \eta\right]\left|\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right|^{-\frac{1}{2}} \frac{\mid 2 \pi \widehat{e}^{\tau_{\beta}}}{}\right|^{-\frac{1}{2}} \\
& \prod_{j} k_{j}! \\
& \times \exp \left(\boldsymbol{k}^{\dagger} \beta(\boldsymbol{z})-\boldsymbol{\rho}^{\dagger} e^{\beta(\boldsymbol{z})}-\frac{1}{2} \beta^{\dagger} \mathcal{F}^{\dagger} e^{-\tau_{\beta}} \mathcal{F} \beta\right. \\
&-\frac{1}{2} \boldsymbol{y}^{\dagger}\left(\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right)^{-1} \boldsymbol{y}+\frac{1}{2} \mathbf{1}^{\dagger} \eta(\boldsymbol{x})-  \tag{3.67}\\
&\left.-\frac{1}{2 \sigma_{\beta}} \tau_{\beta}^{\dagger} \Delta^{\dagger} \Delta \tau_{\beta}-\frac{1}{2 \sigma_{f}} \tau_{f}^{\dagger} \Delta^{\dagger} \Delta \tau_{f}-\frac{1}{2 \sigma_{\eta}} \eta \nabla^{\dagger} \nabla \eta\right)
\end{align*}
$$



Figure 3.3.: Overview over the Bayesian hierarchical model for the case $X \rightarrow Y$. Here, besides the model itself, only strength parameters for the smoothness of the fields are fixed as Hyperparameters $\left(H_{1}\right)$.

### 3.4.5. Simultaneous Laplace Approximation

## The energy

Once more we want to tackle the integration of Eq. 3.67 via a Laplace approximation. In this case however, we cannot divide the integration over $\beta, \tau_{\beta}$ into two separate ones, same as for the integration over $\eta, \tau_{f}$. We will therefore simultaneously find the arguments which maximize the expression under the integral. To do so we rewrite the above expression into an exponential form under the integrals, separating the integrations over $\beta, \tau_{\beta}$ and $\eta, \tau_{f}$ :

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{d} \mid H_{1}\right)= & \frac{1}{\prod_{j} k_{j}!}\left[\int \mathcal { D } [ \beta , \tau _ { \beta } ] \operatorname { e x p } \left(-\frac{1}{2} \log \left|2 \pi e^{\widehat{\tau_{\beta}}}\right|+\boldsymbol{k}^{\dagger} \beta(\boldsymbol{z})-\boldsymbol{\rho}^{\dagger} e^{\beta(\boldsymbol{z})}\right.\right. \\
& \left.\left.-\frac{1}{2} \beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \mathcal{F} \beta}-\frac{1}{2 \sigma_{\beta}} \tau_{\beta}^{\dagger} \Delta^{\dagger} \Delta \tau_{\beta}\right)\right]\left[\int \mathcal { D } [ \eta , \tau _ { f } ] \operatorname { e x p } \left(-\frac{1}{2} \log \left|\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right|\right.\right. \\
& -\frac{1}{2} \boldsymbol{y}^{\dagger}\left(\tilde{F}\left[\tau_{f}\right]+\widehat{\left.\left.\left.e^{\eta(\boldsymbol{x})}\right)^{-1} \boldsymbol{y}+\frac{1}{2} \mathbf{1}^{\dagger} \eta(\boldsymbol{x})-\frac{1}{2 \sigma_{f}} \tau_{f}^{\dagger} \Delta^{\dagger} \Delta \tau_{f}-\frac{1}{2 \sigma_{\eta}} \eta \nabla^{\dagger} \nabla \eta\right)\right] \equiv}\right. \\
& \equiv \frac{1}{\prod_{j} k_{j}!}\left[\int \mathcal{D}\left[\beta, \tau_{\beta}\right] e^{-\gamma_{\xi}\left[\beta, \tau_{\beta}\right]}\right]\left[\int \mathcal{D}\left[\eta, \tau_{f}\right] e^{-\gamma_{\zeta}\left[\eta, \tau_{f}\right]}\right] \tag{3.68}
\end{align*}
$$

Above we exploited that the probability density under the integrals factorizes into terms depending on either $\beta, \tau_{\beta}$ or $\eta, \tau_{f}$, i.e. either on fields that determine the cause distribution or the causal mechanism, but not on both. We write $\mathbb{R}^{[0,1]} \oplus \mathbb{R}^{\mathbb{R}+} \ni \xi \equiv\left(\beta, \tau_{\beta}\right)$ and $\mathbb{R}^{[0,1]} \oplus \mathbb{R}^{\mathbb{R}_{+}} \ni \zeta \equiv\left(\eta, \tau_{f}\right)$.

We introduced the energy functionals:

$$
\begin{align*}
\gamma_{\xi}\left[\beta, \tau_{\beta}\right] \equiv & \frac{1}{2} \log \left|2 \pi e^{\widehat{\tau_{\beta}}}\right|-\boldsymbol{k}^{\dagger} \beta(\boldsymbol{z})+\boldsymbol{\rho}^{\dagger} e^{\beta(\boldsymbol{z})}+\frac{1}{2} \beta^{\dagger} \mathcal{F}^{\dagger} e^{-\tau_{\beta}} \mathcal{F} \beta+\frac{1}{2 \sigma_{\beta}} \tau_{\beta}^{\dagger} \Delta^{\dagger} \Delta \tau_{\beta}  \tag{3.69}\\
\gamma_{\zeta}\left[\eta, \tau_{f}\right] \equiv & \frac{1}{2} \log \left|\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right|+\frac{1}{2} \boldsymbol{y}^{\dagger}\left(\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right)^{-1} \boldsymbol{y}-\frac{1}{2} \mathbf{1}^{\dagger} \eta(\boldsymbol{x}) \\
& +\frac{1}{2 \sigma_{f}} \tau_{f}^{\dagger} \Delta^{\dagger} \Delta \tau_{f}+\frac{1}{2 \sigma_{\eta}} \eta \nabla^{\dagger} \nabla \eta \tag{3.70}
\end{align*}
$$

We consider the fields that minimize the energies:

$$
\begin{align*}
& \xi_{0} \equiv \underset{\xi \in \mathbb{R}^{[0,1]} \not \mathbb{R}^{\mathbb{R}_{+}}}{\operatorname{argmin}} \gamma_{\xi}[\xi]  \tag{3.71}\\
& \zeta_{0} \equiv \underset{\zeta \in \mathbb{R}^{[0,1]} \oplus \mathbb{R}^{\mathbb{R}_{+}}}{\operatorname{argmin}} \gamma_{\zeta}[\zeta] \tag{3.72}
\end{align*}
$$

Again we perform a second-order expansion of $\gamma_{\xi}, \gamma_{\zeta}$ around $\xi_{0}, \zeta_{0}$ w.r.t. $\xi$ and $\zeta$ respectively. In the following we will denote the curvatures by $\Gamma_{\xi}[\xi] \equiv \frac{\delta^{2}}{\delta \xi^{\dagger} \delta \xi} \gamma_{\xi}[\xi]$ and $\Gamma_{\zeta}[\zeta] \equiv \frac{\delta^{2}}{\delta \zeta^{\dagger} \delta \zeta} \gamma_{\zeta}[\zeta]$.

Omitting terms of higher order and exploiting the vanishing gradient (and thus first order) we
end up with:

$$
\begin{align*}
\mathcal{P}\left(\boldsymbol{d} \mid H_{1}\right)= & \frac{1}{\prod_{j} k_{j}!}\left[\int \mathcal{D}[\xi] e^{-\gamma_{\xi}\left[\xi_{0}\right]-\left(\left.\frac{\delta}{\delta_{\xi}} \gamma_{\xi}[\xi]\right|_{\xi=\xi_{0}}\right) \xi-\left.\frac{1}{2} \xi^{\dagger} \Gamma_{\xi}[\xi]\right|_{\xi=\xi_{0}} \xi+\mathcal{O}\left(\xi^{3}\right)}\right] \times \\
& \times\left[\int \mathcal{D}[\zeta] e^{\left.-\gamma_{\zeta}\left[\zeta_{0}\right]-\left(\left.\frac{\delta}{\delta_{\zeta}} \gamma_{\zeta}[\zeta]\right|_{\zeta=\zeta_{0}}\right) \zeta_{-\frac{1}{2} \zeta^{\dagger} \Gamma_{\zeta}[\zeta]| |_{\zeta=\zeta_{0}} \zeta_{+} \mathcal{O}\left(\zeta^{3}\right)}\right] \approx}\right. \\
\approx & \left.\left.\frac{1}{\prod_{j} k_{j}!}\left|\frac{1}{2 \pi} \Gamma_{\xi}[\xi]\right|_{\xi_{0}}\right|^{-\frac{1}{2}}\left|\frac{1}{2 \pi} \Gamma_{\zeta}[\zeta]\right|_{\zeta_{0}}\right|^{-\frac{1}{2}} e^{-\gamma_{\xi}\left[\xi_{0}\right]-\gamma_{\zeta}\left[\zeta_{0}\right]} \tag{3.73}
\end{align*}
$$

To specify the approximation we need to compute the derivatives of first and second order

## Derivatives of $\gamma_{\xi}$

We give the explicit computations in A.1, stating only the results here. We have the gradient:

$$
\begin{array}{r}
\partial_{\beta_{u}} \gamma_{\xi}\left[\beta, \tau_{\beta}\right]=-\boldsymbol{k}_{u}+\rho\left(e^{\beta(\boldsymbol{z})}\right)^{\dagger}+\left(\beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \mathcal{F}}\right)_{u} \\
\partial_{\tau_{\beta u}} \gamma_{\xi}\left[\beta, \tau_{\beta}\right]=\frac{1}{2}-\frac{1}{2} \beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta} \delta_{u}} \mathcal{F} \beta+\left(\frac{1}{\sigma_{\beta}} \tau_{\beta}^{\dagger} \Delta^{\dagger} \Delta\right)_{u}} \tag{3.75}
\end{array}
$$

The curvature, $\Gamma_{\xi}\left[\beta, \tau_{\beta}\right]$ has non-vanishing mixed derivatives,

$$
\Gamma_{\xi}\left[\beta, \tau_{\beta}\right]=\left(\begin{array}{cc}
\frac{\delta^{2} \gamma_{\xi}}{\delta \beta^{\dagger} \delta \beta} & \frac{\delta^{2} \gamma_{\xi}}{\delta \beta^{\dagger} \tau_{\beta}}  \tag{3.76}\\
\frac{\delta^{2} \gamma_{\xi}}{\delta \tau_{\beta}^{\dagger} \delta \beta} & \frac{\delta^{2} \tau_{\xi}}{\delta \tau_{\beta}^{\dagger} \delta \tau_{\beta}}
\end{array}\right)
$$

We have the terms

$$
\begin{align*}
\partial_{\beta_{u}} \partial_{\beta_{v}} \gamma_{\xi}\left[\beta, \tau_{\beta}\right] & =\left(\widehat{\rho e^{\beta(\boldsymbol{z})}}\right)_{u v}+\left(\mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}}} \mathcal{F}\right)_{u v}  \tag{3.77}\\
\partial_{\tau_{\beta_{u}}} \partial_{\tau_{\beta_{v}}} \gamma_{\xi}\left[\beta, \tau_{\beta}\right] & =\frac{1}{2} \delta_{u v} \beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \delta_{z}} \mathcal{F} \beta+\left(\frac{1}{\sigma_{\beta}} \Delta^{\dagger} \Delta\right)_{u v}  \tag{3.78}\\
\partial_{\tau_{\beta u}} \partial_{\beta_{v}} \gamma_{\xi}\left[\beta, \tau_{\beta}\right] & =-\left(\beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta} \delta_{v}} \mathcal{F}}\right)_{u}=\partial_{\beta_{v}} \partial_{\tau_{\beta_{u}}} \gamma_{\xi}\left[\beta, \tau_{\beta}\right] \tag{3.79}
\end{align*}
$$

Exploiting that for some block matrix we have $(\| \operatorname{Sil00})$ :

$$
\operatorname{det}\left(\begin{array}{ll}
A & B  \tag{3.80}\\
C & D
\end{array}\right)=\operatorname{det}(A) \operatorname{det}\left(D-C A^{-1} B\right)=\operatorname{det}(D) \operatorname{det}\left(A-B D^{-1} C\right)
$$

under the condition that both, $A$ and $D$, are invertible. we get:

$$
\begin{equation*}
\operatorname{det} \Gamma_{\xi}=\left(\left|\frac{\delta^{2} \gamma}{\delta \beta^{\dagger} \delta \beta}\right|\left|\frac{\delta^{2} \gamma}{\delta \tau_{\beta}^{\dagger} \delta \tau_{\beta}}-\frac{\delta^{2} \gamma}{\delta \tau_{\beta}^{\dagger} \delta \beta}\left(\frac{\delta^{2} \gamma}{\delta \beta^{\dagger} \delta \beta}\right)^{-1} \frac{\delta^{2} \gamma}{\delta \beta^{\dagger} \delta \tau_{\beta}}\right|\right) \tag{3.81}
\end{equation*}
$$

## Derivatives of $\gamma_{\zeta}$

Making the definitions

$$
\begin{align*}
& \Lambda: \mathbb{R} \rightarrow \mathbb{C}^{N \times N} \\
& \Lambda(u)_{i j}=\left(\mathcal{F}^{\dagger} \widehat{e^{\tau_{f}} \delta_{u}} \mathcal{F}\right)_{x_{i} x_{j}} \tag{3.82}
\end{align*}
$$

Again we refer to the explicit computation of the derivatives in A.2, only giving the results here For the gradient we get:

$$
\left.\begin{array}{rl}
\partial_{\eta_{u}} \gamma_{\zeta}\left[\eta, \tau_{f}\right] & =\frac{1}{2} \operatorname{tr}\left(G e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}\right.
\end{array}\right)-\frac{1}{2} \boldsymbol{y}^{\dagger} G e^{\eta(\boldsymbol{x})} G \delta_{\boldsymbol{x} u} \boldsymbol{y}-\frac{1}{2}(\boldsymbol{x})+\left(\frac{1}{\left.\sigma_{\eta} \eta^{\dagger} \nabla^{\dagger} \nabla\right)_{u}} \begin{array}{r}
\partial_{\tau_{f_{u}}} \gamma_{\zeta}\left[\eta, \tau_{f}\right] \\
=\frac{1}{2} \operatorname{tr}\left(G \Lambda_{u}\right)-\frac{1}{2} \boldsymbol{y}^{\dagger} G \Lambda_{u} G \boldsymbol{y}+\left(\frac{1}{\sigma_{f}} \tau_{f}^{\dagger} \Delta^{\dagger} \Delta\right)_{u} \tag{3.84}
\end{array}\right.
$$

And for the curvature,

$$
\Gamma_{\zeta}\left[\eta, \tau_{f}\right]=\left(\begin{array}{cc}
\frac{\delta^{2} \gamma_{\zeta}}{\delta \eta^{\top} \dagger} \frac{\delta^{2} \gamma_{\zeta}}{\delta \eta} & \frac{\delta^{\prime}}{\delta \eta^{\dagger} \delta \tau_{f}} \\
\frac{\delta^{2} \gamma_{\zeta}}{\delta \tau_{f}^{\dagger} \delta \eta} & \frac{\delta^{2} \zeta_{\zeta}}{\delta \tau_{f}^{\top} \delta \tau_{f}}
\end{array}\right)
$$

We have the terms

$$
\begin{align*}
& \partial_{\eta_{u}} \partial_{\eta_{v}} \gamma_{\zeta}\left[\eta, \tau_{f}\right]=\frac{1}{2} \operatorname{tr}\left(-G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}}} G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}} v}+\delta_{u v} G e^{\widehat{\eta(\boldsymbol{x}) \delta_{\boldsymbol{x}}}}\right) \\
& +\frac{1}{2} \boldsymbol{y}^{\dagger}\left(2 G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}} u} G e^{\hat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} v}} G-G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}}} \delta_{u v} G\right) \boldsymbol{y}+\left(\frac{1}{\sigma_{\eta}} \nabla^{\dagger} \nabla\right)_{u v}  \tag{3.86}\\
& \partial_{\tau_{f_{u}}} \partial_{\tau_{f_{v}}} \gamma_{\zeta}\left[\eta, \tau_{f}\right]=\frac{1}{2} \operatorname{tr}\left(-G \Lambda_{u} G \Lambda_{v}+G \Lambda_{u} \delta_{u v}\right) \\
& +\frac{1}{2} \boldsymbol{y}^{\dagger}\left(2 G \Lambda_{u} G \Lambda_{v} G-G \Lambda_{u} \delta_{u v} G\right) \boldsymbol{y}+\left(\frac{1}{\sigma_{f}} \Delta^{\dagger} \Delta\right)_{u v}  \tag{3.87}\\
& \partial_{\eta_{u}} \partial_{\tau_{f_{v}}} \gamma_{\zeta}\left[\eta, \tau_{f}\right]=-\frac{1}{2} \operatorname{tr}\left(G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G \Lambda_{v}\right)+\boldsymbol{y}^{\dagger} G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G \Lambda_{v} G \boldsymbol{y}=\partial_{\tau_{f_{v}}} \partial_{\eta_{u}} \gamma_{\zeta}\left[\eta, \tau_{f}\right] \tag{3.88}
\end{align*}
$$

Using again 3.80, we get:

$$
\begin{equation*}
\operatorname{det} \Gamma_{\zeta}=\left(\left|\frac{\delta^{2} \gamma}{\delta \tau_{f}^{\dagger} \delta \tau_{f}}\right|\left|\frac{\delta^{2} \gamma}{\delta \eta^{\dagger} \delta \eta}-\frac{\delta^{2} \gamma}{\delta \eta^{\dagger} \delta \tau_{f}}\left(\frac{\delta^{2} \gamma}{\delta \tau_{f}^{\dagger} \delta \tau_{f}}\right)^{-1} \frac{\delta^{2} \gamma}{\delta \tau_{f}^{\dagger} \delta \eta}\right|\right) \tag{3.89}
\end{equation*}
$$

We do not refrain the explicit terms at this place and refer to A. 2 where these are explicitly given.

### 3.4.6. Bayes Factor for the Deep Model Selection

Plugging in the approximations above we can state the information Hamiltonian for the evidence in a deep model selection, i.e. the only fixed hyperparameters are strength parameters for the smoothness enforcing priors $\sigma_{\beta}, \sigma_{f}, \sigma_{\eta}$ :

$$
\begin{align*}
\mathcal{H}\left(\boldsymbol{d} \mid H_{1}\right) \approx & \left.\mathcal{H}_{0}+\log \prod_{j} k_{j}!+\frac{1}{2} \log \left|\frac{1}{2 \pi} \Gamma_{\xi}[\xi]\right|_{\xi_{0}} \right\rvert\,+ \\
& \left.+\frac{1}{2} \log \left|\frac{1}{2 \pi} \Gamma_{\zeta}[\zeta]\right| \zeta_{0} \right\rvert\,+\gamma_{\xi}\left[\xi_{0}\right]+\gamma_{\zeta}\left[\zeta_{0}\right] \tag{3.90}
\end{align*}
$$

The model evidence for the reverse direction is again obtained by switching $\boldsymbol{x}$ and $\boldsymbol{y}$ in all terms. The Bayes factor is then given by:

$$
\begin{equation*}
\mathcal{O}_{X \rightarrow Y}=\exp \left[-\mathcal{H}\left(\boldsymbol{d} \mid H_{1}\right)+\mathcal{H}\left(\boldsymbol{d} \mid H_{-1}\right)\right] \tag{3.91}
\end{equation*}
$$

If $\mathcal{O}_{X \rightarrow Y}>1$, the Bayesian Model Selection suggests the causal direction $X \rightarrow Y$, if $\mathcal{O}_{X \rightarrow Y}<1$, the other direction, $Y \rightarrow X$ is suggested instead. In case of equality, $\mathcal{O}_{X \rightarrow Y}=1$, no direction is preferred. If one admits independence of the variables, i.e. neither $X \rightarrow Y$ or $Y \rightarrow X$, the Bayes factor being equal to 1 or very close would an indicator for this latter case.

## 4. Implementation and Benchmarks

### 4.1. Sampling Causal Data via a Forward Model

To estimate the performance of our algorithm and compare it with other existing approaches, a benchmark dataset is of interest to us. Such benchmark data is usually either real world data, or synthetically produced one. While we will use the TCEP benchmark set Moo+16 in 4.4.5, we also want to use our outlined formalism to generate artificial data representing causal structures. Based on our derivation for cause and effect we implement a forward model to generate data $\boldsymbol{d}$ as following:

Algorithm 1. Sampling of causal data via forward model
Input: Power spectra $P_{\beta}, P_{f}$,
noise variance $\varsigma^{2}$, number of bins $n_{\text {bins }}$,
approximate desired number of samples $\tilde{N}$
Output: $N$ samples $\left(d_{i}\right)=\left(x_{i}, y_{i}\right)$ generated from a causal relation of either $X \rightarrow Y$ or $Y \rightarrow X$

1. Draw a sample field $\beta \in \mathbb{R}^{[0,1]}$ from the distribution $\mathcal{N}(\beta \mid 0, B)$
2. Set an equally spaced grid with $n_{\text {bins }}$ points in the interval $[0,1]: \boldsymbol{z}=\left(z_{1}, \ldots, z_{n_{\text {bins }}}\right), z_{i}=$ $\frac{i-0.5}{n_{\text {bins }}}$
3. Calculate the vector of Poisson means $\boldsymbol{\lambda}=\left(\lambda_{1}, \ldots \lambda_{n_{\text {bins }}}\right)$ with $\lambda_{i} \propto e^{\beta\left(z_{i}\right)}$
4. At each grid point $i \in\left\{1, \ldots, n_{\text {bins }}\right\}$, draw a sample $k_{i}$ from a Poisson distribution with mean $\lambda_{i}: k_{i} \sim \mathcal{P}_{\lambda_{i}}\left(k_{i}\right)$
5. Set $N=\sum_{i=1}^{n_{\text {bins }}} k_{i}$,
6. For each $i \in\left\{1, \ldots, n_{\text {bins }}\right\}$ add $k_{i}$ times the element $z_{i}$ to the set of measured $x_{j}$. Construct the vector $\boldsymbol{x}=(\ldots, \underbrace{z_{i}, z_{i}, z_{i}}_{k_{i} \text { times }}, \ldots)$
7. Draw a sample field $f \in \mathbb{R}^{[0,1]}$ from the distribution $\mathcal{N}(f \mid 0, F)$. Rescale $f$ s.th. $f \in[0,1]^{[0,1]}$.
8. Draw a multivariate noise sample $\boldsymbol{\epsilon} \in \mathbb{R}^{N}$ from a normal distribution with zero mean

$$
\text { and variance } \varsigma^{2}, \boldsymbol{\epsilon} \sim \mathcal{N}\left(\boldsymbol{\epsilon} \mid 0, \varsigma^{2}\right)
$$

9. Generate the effect data $\boldsymbol{y}$ by applying $f$ to $\boldsymbol{x}$ and adding $\boldsymbol{\epsilon}: \boldsymbol{y}=f(\boldsymbol{x})+\boldsymbol{\epsilon})$
10. With probability $\frac{1}{2}$ return $\boldsymbol{d}=\left(\boldsymbol{x}^{T}, \boldsymbol{y}^{T}\right)$, otherwise return $\boldsymbol{d}=\left(\boldsymbol{y}^{T}, \boldsymbol{x}^{T}\right)$,
[^3]Comparing the samples for different power spectra (see Fig. 4.1), we decide to sample data with power spectra $P(q)=\frac{1000}{q^{4}+1}$ and $P(q)=\frac{1000}{q^{6}+1}$, as these seem to resemble "natural" mechanisms.

### 4.2. Implementation of the Bayesian Causal Inference

### 4.2.1. Implementation of the Shallow Model

Based on our derivation in 3.2 we propose a specific algorithm to decide the causal direction of a given dataset and therefore give detailed answer for problem 1. Basically the task comes down to find the minimum $\beta_{0}$ for the saddle point approximation and calculate the terms given in Eq. 3.36

Algorithm 2. 2-variable causal inference
Input: Finite sample data $\boldsymbol{d} \equiv(\boldsymbol{x}, \boldsymbol{y}) \in \mathbb{R}^{N \times 2}$, Hyperparameters $P_{\beta}, P_{f}, \varsigma^{2}, r$
Output: Predicted causal direction $\mathcal{D}_{X \rightarrow Y} \in\{-1,1\}$ where -1 represents the prediction " $Y \rightarrow X$ " and 1 represents $X \rightarrow Y$

1. Rescale the data to the $[0,1]$ interval. I.e. $\min \left\{x_{1}, \ldots, x_{N}\right\}=\min \left\{y_{1}, \ldots, y_{N}\right\}=0$ and $\max \left\{x_{1}, \ldots, x_{N}\right\}=\max \left\{y_{1}, \ldots, y_{N}\right\}=1$
2. Define an equally spaced grid of $\left(z_{1}, \ldots, z_{n_{\text {bins }}}\right)$ in the interval $[0,1]$
3. Calculate matrices $\boldsymbol{B}, \boldsymbol{F}$ representing the covariance operators $B$ and $F$ evaluated at the positions of the grid, i.e. $\boldsymbol{B}_{i j}=B\left(z_{i}, z_{j}\right)$
4. Find the $\beta_{0} \in \mathbb{R}^{[0,1]}$ s.th. $\gamma$ in Eq. 3.15 becomes minimal
5. Calculate the $\boldsymbol{d}$-dependent terms of the information Hamiltonian in Eq. 3.36 (i.e. all besides $\mathcal{H}_{0}$ )
6. Repeat steps 4 and 5 with $\boldsymbol{y}$ and $\boldsymbol{x}$ switched

Figure 4.1.: Different sampled fields from $\mathcal{N}\left(\cdot \mid 0, \mathcal{F}^{\dagger} P \mathcal{F}\right)$ with the power spectrum $P(q) \propto \frac{1}{q^{2}+1}$ (top), $P(q) \propto \frac{1}{q^{4}+1}$ (middle), $P(q) \propto \frac{1}{q^{6}+1}$ (bottom). On the left, the field values themselves are plotted, on the right an exponential function is applied as in our formulation for $\lambda_{j} \propto e^{\beta\left(z_{j}\right)}$

(a) fields sampled with $P \propto \frac{1}{q^{2}+1}$

(c) fields sampled with $P \propto \frac{1}{q^{4}+1}$

(e) fields sampled with $P \propto \frac{1}{q^{6}+1}$

(b) exponential function applied to fields sampled with $P \propto \frac{1}{q^{2}+1}$

(d) exponential function applied to fields sampled with $P \propto \frac{1}{q^{4}+1}$

(f) exponential function applied to fields sampled with $P \propto \frac{1}{q^{6}+1}$

Figure 4.3.: Synthetic Datasets sampled via alg. 1. Blue scatter plots indicate a true causal direction of $X \rightarrow Y$, red scatter plots the direction $Y \rightarrow X$

7. Calculate the Bayes factor $\mathcal{O}_{X \rightarrow Y}$
8. If $\mathcal{O}_{X \rightarrow Y}>1$, return 1, else return -1

We provide an implementation of alg. 2 in Python ${ }^{2}$. We approximate the operators $B, F$ as matrices $\in \mathbb{R}^{n_{\text {bins }} \times n_{\text {bins }}}$, which allows us to explicitly numerically compute the determinants and the inverse. As the most critical part we consider the minimization of $\beta$, i.e. step 4 in 2 . As we are however able to analytically give the curvature $\Gamma_{\beta}$ and the gradient $\partial_{\beta} \gamma_{\beta}$, we can use a Newton-scheme here. We derive satisfying results (see Fig. 4.4) using the Newton-CG algorithm NW06, provided by the SciPy-Library [J+01]. After testing our algorithm on different benchmark data, we choose the default hyperparameters as

$$
\begin{align*}
& P_{\beta}=P_{f} \propto \frac{1}{q^{4}+1},  \tag{4.1}\\
& \varsigma^{2}=0.01,  \tag{4.2}\\
& r=512,  \tag{4.3}\\
& \rho=1 . \tag{4.4}
\end{align*}
$$

While fixing the power spectra might seem somewhat arbitrary, we remark that this corresponds to fixing a kernel e.g. as a squared exponential kernel, which is done in many publications (e.g. MST18, Gou+17])

### 4.2.2. Issues with the Implementation of the Deeper Models

We tested the possibility to implement the deeper models. In the noise inference model, as described in 3.3 , we need to go through two separate minimization phases. While the first minimization is the same one as 4 in 2, the second one, i.e. determining $\eta_{0}$, involves the explicit numerical inversion of matrices $\in \mathbb{R}^{n_{\text {bins }} \times n_{\text {bins }}}$ which depend on $\eta$ when calculating the curvature $\Gamma_{\eta}[\eta]$. Using again a Newton-scheme which involves several computations of the curvature in each minimization step, this minimization becomes therefore vastly computationally expensive and such too slow to handle it efficiently.

Implementing the "deep" model as outlined in 3.4 would now involve finding numerical representations for $\beta, \tau_{\beta}, \eta, \tau_{f}$. Besides the computational complexity which would behave even worse than mentioned above, one now has the difficulty that the minima $\beta_{0}$ and $\tau_{\beta}$ depend on each other, i.e. a change in $\tau_{\beta}$ will lead to a different $\beta_{0}$. The same holds for $\eta$ and $\tau_{f}$. We therefore leave these models at a theoretical state here and consider their implementation as an option for future work to be done.

[^4]
(a) Scatter plot for (synthetic) causal data, the (b) histogram ( $\boldsymbol{k}$ ) and result of the computatrue direction is $X \rightarrow Y$

(c) histogram ( $\boldsymbol{k}$ ) and result of the computa-(d)
tional minimization of $\beta$ for the model in the direction $X \rightarrow Y$

 histogram $(\boldsymbol{k})$ and result of the computa-
tional minimization of $\beta$ for the model in the direction $Y \rightarrow X$

Figure 4.4.: Demonstration of the computational minimization

### 4.3. Methods for Comparison

We compare our outlined model, in the following called BCM, short for Bayesian Causal Model, to a number of state-of-the-art approaches. The selection of the considered methods is influenced by the ones in recent publications, e.g. MST18; Gou+17.

As it is one of the oldest algorithms in the field and has been used as comparison in a variety of publications, we include the LiNGAM algorithm. We also use the ANM algorithm [Moo+16] with HSIC and Gaussian Process Regression (ANM-HSIC)

The ANM-MML approach [Ste+10] uses a Bayesian Model Selection, from the perspective of formulation it is the closest to the algorithm proposed within this thesis, at least to our best knowledge. This makes it an interesting reference and motivates the choice to include it here.

We further include the IGCI algorithm, as it differs fundamentally in its formulation from the ANM algorithms and has shown strong results in recent publications Moo+16; Gou+17; MST18. We employ the IGCI algorithm with entropy estimation for scoring and a Gaussian distribution as reference distribution.

Finally, CGNN Gou+17] represents the rather novel influence of deep learning methods. As it proved to perform generally well in different scenarios Gou+17; MST18], we include it in our comparison. We use the implementation provided by the authors, with itself uses Python with the Tensorflow Aba+15 library. The most critical hyper-parameter here is, as the authors themselves mention, the number of hidden neurons. We set this number to a value of $n_{h}=30$, as this is the default in the given implementation and delivers generally good results. We use 8 runs each, in our eyes this represents a adequate trade-off between unnecessary high computation time and bad performance for reasons of being too restrictive.

A comparison with the KCDC algorithm would be interesting, unfortunately the authors did not provide any computational implementation so far (September 2018).

### 4.4. Results

### 4.4.1. Default Synthetic Data

We arrive at the conclusion to choose spectra of the type $P(q)=\frac{1}{q^{4}+1}$, for both spectra. We further set $n_{\text {bins }}=512, \tilde{N}=300$ and $\varsigma^{2}=0.05$ as default settings for producing synthetic causal data. The scatter plots for the resulting data are shown in 4.3 .

The resulting accuracies are given below:

Table 4.1.: Accuracy for the default synthetic data benchmark.

| Model | accuracy |
| :--- | ---: |
| BCM | 0.98 |
| LiNGAM | 0.30 |
| ANM-HSIC | 1.00 |
| ANM-MML | 1.00 |
| IGCI | 0.65 |
| CGNN | 0.72 |

### 4.4.2. High Noise Data

As a first variation, we explore the influence of high and very high noise on the performance of the inference models. Therefore we set the parameter $\varsigma^{2}=0.2$ for high noise and $\varsigma^{2}=1$ for very high noise in 1, while keeping the other parameters set to the same values as in 4.4.1. The scatter plots for the resulting datasets are given in figures C.1 and fig C.2.

While our BCM algorithm is able to still perform somewhat reliable with a accuracy of $\geq 90 \%$, especially the ANM algorithms do not remarkably suffer from the noise. This is likely due to the fact that the distribution of the true cause $\mathcal{P}(X)$ is not influenced by the high noise and this distribution is assessed in its own.

Table 4.2.: Accuracy for the high-noise synthetic data benchmark.

| Model | $\varsigma^{2}=0.2$ | $\varsigma^{2}=1$ |
| :--- | ---: | ---: |
| BCM | 0.94 | 0.90 |
| LiNGAM | 0.31 | 0.40 |
| ANM-HSIC | 0.98 | 0.94 |
| ANM-MML | 0.99 | 0.99 |
| IGCI | 0.60 | 0.58 |
| CGNN | 0.75 | 0.77 |

### 4.4.3. Strongly Discretized Data

As our model uses a Poissonian approach, which explicitly considers discretization effects of data measurement, it is of interest how the performance behaves when using a strongly discretization. We emulate such a situation by employing our forward model 1 with a very low number of bins. We keep all parameters as in 4.4.1 and set $n_{\text {bins }}=16$ and $n_{\text {bins }}=8$ for synthetic data with high and very high discretization. The visualization of the datasets is given in figures C. 3 and C.4. Apparently, especially the ANM models do not suffer anyhow from the strong discretization. CGNN and IGCI perform significantly worse here. In the case of IGCI this can be explained by the entropy estimation, which simply removes non-unique samples. Our BCM algorithm is able to achieve over $90 \%$ accuracy here.

Table 4.3.: Accuracy for the strongly discretized synthetic benchmark data

| Model | $n_{\text {bins }}=16$ | $n_{\text {bins }}=8$ |
| :--- | ---: | ---: |
| BCM | 0.93 | 0.97 |
| LiNGAM | 0.23 | 0.21 |
| ANM-HSIC | 0.99 | 1.00 |
| ANM-MML | 1.00 | 1.00 |
| IGCI | 0.24 | 0.09 |
| CGNN | 0.57 | 0.22 |

### 4.4.4. Very Sparse Data

We explore another challenge for inference algorithms where we strongly reduce the number of samples. While we sampled about 300 observations with our other forward models so far, here we reduce the number of observed samples to 30 and 10 samples. Again we refer to the scatter plots in figures C.5 and C. 6 . In this case our model performs very well compared to the other models, in fact it is able to outperform them in the case of just 10 samples being given.

We note that of course our model does have the advantage that it "knows" the hyperparameters of the underlying forward model. Yet we consider the results as encouraging.

Table 4.4.: Accuracy for very sparse data.

| Model | 30 samples | 10 samples |
| :--- | ---: | ---: |
| BCM | 0.92 | 0.75 |
| LiNGAM | 0.44 | 0.45 |
| ANM-HSIC | 0.91 | 0.71 |
| ANM-MML | 0.98 | 0.69 |
| IGCI | 0.48 | 0.40 |
| CGNN | 0.46 | 0.39 |

### 4.4.5. Real World Data

The most widely used benchmark set with real world data is the Tübingen Cause Effect Pairs dataset (TCEP) Moo+16. The collection currently (August 2018) consists of 108 datasets. However these include sets with multiple cause or effect variables. Excluding these reduced the collection to 102 datasets. As proposed by the maintainers, we use a weighted evaluation of accuracy here. As some of the datasets represent essentially the same mechanism and just have been collected with different means (e.g. "latitude" and "life expectancy at birth for different countries, female, 1995-2000" in pair0057 vs. "latitude" and life expectancy at birth for different countries, female, 1990-1995" in pair0058) they are assigned a reduced weight. A full description of the TCEP benchmark set is given in D accompanied by scatter plots of

As some of the contained datasets include a high number of samples (up to 11000) we randomly subsample large datasets to 500 samples each in order to keep computation time maintainable. We did not include the LiNGAM algorithm here, as we experienced computational problems with obtaining results here for certain datasets (namely pair0098). (Gou+17] report the accuracy of LiNGAM on the TCEP dataset at around $40 \%$. Our model shows to perform generally comparable to established approaches as $A N M$ and $I G C I . C G N N$ performs best here, the accuracy which we evaluate at about $70 \%$ is however lower than the one reported by the authors Gou +17 at around $80 \%$. The reason for this is arguably to be found in the fact that we set all hyperparameters to fixed values, while Gou+17 used a leave-one-out-approach to find the best setting for the hyperparameter $n_{h}$.

Motivated by the generally strong performance of our approach in the case of sparse data, we also explore a situation where real world data is only sparsely available. To that end, we subsample the TCEP datasets s.th. each 20 randomly chosen samples are kept. We give visual impression for this data in Fig. 4.5 The results are as well given in Table 4.5 The loss in accuracy of our model is remarkably small. In fact, $B C M$ is able to outperform the other models considered here, even if not by a large margin.

Table 4.5.: Accuracy for TCEP Benchmark

| Model | TCEP | TCEP with 20 samples |
| :--- | ---: | ---: |
| BCM | 0.64 | 0.60 |
| ANM-HSIC | 0.63 | 0.41 |
| ANM-MML | 0.58 | 0.51 |
| IGCI | 0.66 | 0.59 |
| CGNN | 0.70 | 0.55 |

Figure 4.5.: Real World Data TCEP benchmark set with 20 randomly chosen samples per dataset. Blue scatter plots indicate a true causal direction of $X \rightarrow Y$, red scatter plots the direction $Y \rightarrow X$


## 5. Discussion and Conclusion

In this thesis we introduced a model for the 2-variable causal inference task. Our model builds on the formalism of information field theory which explicitly models the connection of finite dimensional measurement data to underlying infinite dimensional structures. In this regard, we employed the concept of Bayesian model selection and made the assumption of additive noise, i.e. $\boldsymbol{x}=f(\boldsymbol{y})+\boldsymbol{\epsilon}$. In contrast to other methods which do so, such as ANM-MML, we do not model the cause distribution by a Gaussian mixture model but by a Poisson Lognormal statistic.

We could show that our model is able to provide reliable classification accuracy in the present causal inference task. Another difference from our model to existing ones is arguably to be found in the choice of the covariance operators. While most other publications use squared exponential kernels for Gaussian process regression, we choose a covariance which is governed by a $\frac{1}{q^{4}+1}$ power spectrum. This leads arguably to a different importance of structure at small scales which is emphasized more strongly by our covariance than in a squared exponential kernel.

As a certain weak point within our model we consider the approximation of the uncomputable path integrals via the Laplace approximation. A thorough investigation of error bounds, e.g. via Maj15 is yet to be carried out. As an alternative, one can think about samplingbased approaches to approximate the integrals. A recent publication $(\mid \overline{\mathrm{Cal}+18})$ introduced a harmonic-mean based sampling approach to approximate high dimensional integrals. Such a technique might be promising in the context of our formalism.

Another outlook is to be seen in the computational implementation of the deeper models discussed in 3.3 and 3.4 . However the feasibility of this is certainly questionable. Especially with the outlined technique of the Laplace approximation the challenge persists in the computational minimization which is numerically highly complex in this case.

Yet, the implementation of our model with fixed noise variance and power spectra was able to deliver competitive results with regard to state-of-the-art methods in the benchmarks. In particular, our method seems to be slightly superior in the low sample regime, probably due to the more appropriate Poisson statistic used. We consider this as an encouraging result for a first work in the context of information field theory-based causal inference.

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## A. Explicit Calculations

## A.1. Derivatives of $\gamma_{\xi}$

## A.1.1. The Term log $\left|2 \pi \widehat{e^{\tau_{\beta}}}\right|$

$$
\begin{align*}
\partial_{\tau_{\beta}} \log \left|2 \pi \widehat{e^{\tau_{\beta}}}\right| & =\left|2 \pi \widehat{e^{\tau_{\beta}}}\right|^{-1} \partial_{\tau_{\beta_{q}}}\left|2 \pi \widehat{e^{\tau_{\beta}}}\right| \\
& =\left|2 \pi \widehat{e^{\tau_{\beta}}}\right|^{-1}\left|2 \pi e^{\widehat{\tau}_{\beta}}\right| \operatorname{tr}\left(\widehat{e^{-\tau_{\beta}}} \partial_{\tau_{\beta_{q}}} \widehat{e^{\tau_{\beta}}}\right) \\
& =\operatorname{tr}\left(\widehat{e^{-\tau_{\beta}}} \widehat{e^{\tau_{\beta} \delta_{q}}}\right) \\
& =\operatorname{tr}\left(\widehat{\delta_{q}}\right)=1 \tag{A.1}
\end{align*}
$$

The second order derivative w.r.t. $\tau_{\beta}$ therefore vanishes,

$$
\begin{equation*}
\partial_{\tau_{\beta_{r}}} \partial_{\tau_{\beta_{q}}} \log \left|2 \pi \widehat{e}^{\widehat{\tau_{\beta}}}\right|=0 \tag{A.2}
\end{equation*}
$$

As well as the second order mixed derivatives,

$$
\begin{equation*}
\partial_{\beta_{u}} \partial_{\tau_{\beta_{q}}} \log \left|2 \pi e^{\widehat{\tau_{\beta}}}\right|=0=\partial_{\tau_{\beta_{q}}} \partial_{\beta_{u}} \log \left|2 \pi e^{\widehat{\tau_{\beta}}}\right| \tag{A.3}
\end{equation*}
$$

## A.1.2. The Term $\beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}}} \mathcal{F} \beta$

As first order derivatives we have

$$
\begin{align*}
\partial_{\beta_{u}} \beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}}} \mathcal{F} \beta & =\left(\beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \mathcal{F}}\right)_{u}  \tag{A.4}\\
\partial_{\tau_{\beta q}} \beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}}} \mathcal{F} \beta & =\beta^{\dagger} \mathcal{F}^{\dagger}\left(\partial_{\tau_{\beta q}} \widehat{e^{-\tau_{\beta}}}\right) \mathcal{F} \beta=-\beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \delta_{q}} \mathcal{F} \beta \tag{A.5}
\end{align*}
$$

The above derivative w.r.t. to $\beta$ is trivial, in the derivative w.r.t. $\tau_{\beta}$ we used that the Fourier transform is invariant under the derivative. The resulting expression $\widehat{e^{-\tau_{\beta}} \delta_{q}}$ is to be understood as the field analogue of a diagonal matrix where only a single entry (with the index $q$ ) is not zero. Therefore the second order derivatives are:

$$
\begin{align*}
\partial_{\beta_{v}} \partial_{\beta_{u}} \beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \mathcal{F}} \beta & =\left(\mathcal{F}^{\dagger} \widehat{\left.e^{-\tau_{\beta}} \mathcal{F}\right)_{u v}}\right.  \tag{A.6}\\
\partial_{\tau_{\beta_{r}}} \partial_{\tau_{\beta q}} \beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \mathcal{F} \beta} & =\delta_{q r} \beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \delta_{q}} \mathcal{F} \beta  \tag{A.7}\\
\partial_{\tau_{\beta_{q}}} \partial_{\beta_{u}} \beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \mathcal{F} \beta} & =\partial_{\tau_{\beta_{q}}}\left(\beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}}} \mathcal{F}\right)_{u}=-\left(\beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \delta_{q}} \mathcal{F}\right)_{u} \tag{A.8}
\end{align*}
$$

## A.2. Derivatives of $\gamma_{\zeta}$

## A.2.1. The Term $\boldsymbol{y}^{\dagger}\left(\tilde{F}\left[\tau_{f}\right]+e^{\widehat{\eta(\boldsymbol{x}})}\right)^{-1} \boldsymbol{y}$

The derivatives of the $\tau_{f}$-dependent matrix elements $\tilde{F}_{i j}$ are

$$
\begin{align*}
\partial_{\tau_{f_{q}}} \tilde{F}\left[\tau_{f}\right]_{i j} & =\partial_{\tau_{f_{q}}}\left(\mathcal{F}^{\dagger} \widehat{e^{\tau_{f}} \mathcal{F}}\right)_{x_{i} x_{j}}=\left(\mathcal{F}^{\dagger} \widehat{e^{\tau_{f}} \delta_{q}} \mathcal{F}\right)_{x_{i} x_{j}} \equiv\left(\Lambda_{q}\right)_{i j}  \tag{A.9}\\
\partial_{\tau_{f_{r}}} \partial_{\tau_{f_{q}}} \tilde{F}\left[\tau_{f}\right]_{i j} & =\partial_{\tau_{f_{r}}}\left(\mathcal{F}^{\dagger} \widehat{e^{\tau_{f} \delta_{q}} \mathcal{F}}\right)_{x_{i} x_{j}}=\left(\Lambda_{q}\right)_{i j} \delta_{q r} \tag{A.10}
\end{align*}
$$

For the sake of brevity we introduced $\left(\Lambda_{q}\right)_{i j} \equiv\left(\mathcal{F}^{\dagger} \widehat{e^{\tau_{f}} \delta_{q}} \mathcal{F}\right)_{x_{i} x_{j}}$ above, which we will use from now on. We will further write $G \equiv\left(\tilde{F}+\widehat{e^{\eta(\boldsymbol{x})}}\right)^{-1}$

$$
\begin{align*}
\partial_{\tau_{f_{q}}} \boldsymbol{y}^{\dagger} G \boldsymbol{y} & =-\boldsymbol{y}^{\dagger} G \Lambda_{q} G \boldsymbol{y}  \tag{A.11}\\
\partial_{\tau_{f_{r}}} \partial_{\tau_{f_{q}}} \boldsymbol{y}^{\dagger} G \boldsymbol{y} & =\boldsymbol{y}^{\dagger}\left(2 G \Lambda_{q} G \Lambda_{r} G-G \Lambda_{q} \delta_{q r} G\right) \boldsymbol{y}  \tag{A.12}\\
\partial_{\eta_{u}} \boldsymbol{y}^{\dagger} G \boldsymbol{y} & =-\boldsymbol{y}^{\dagger} G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G \boldsymbol{y}  \tag{A.13}\\
\partial_{\eta_{v}} \partial_{\eta_{u}} \boldsymbol{y}^{\dagger} G \boldsymbol{y} & =\boldsymbol{y}^{\dagger}\left(2 G \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G \widehat{e^{\eta(\boldsymbol{x})}} \delta_{\boldsymbol{x} r} G-G \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}} u} \delta_{u v} G\right) \boldsymbol{y}  \tag{A.15}\\
\partial_{\eta_{u}} \partial_{\tau_{f_{q}}} \boldsymbol{y}^{\dagger} G \boldsymbol{y} & =-2 \boldsymbol{y}^{\dagger}\left(\partial_{\eta_{u}} G\right) \Lambda_{q} G \boldsymbol{y}=2 \boldsymbol{y}^{\dagger} G \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G \Lambda_{q} G \boldsymbol{y}  \tag{A.16}\\
& =\partial_{\tau_{f_{q}}} \partial_{\eta_{u}} \boldsymbol{y}^{\dagger} G \boldsymbol{y}
\end{align*}
$$

A.2.2. The Term $\log \left|\tilde{F}\left[\tau_{f}\right]+e^{\widehat{\eta(\boldsymbol{x})}}\right|$

$$
\begin{align*}
& \partial_{\tau_{f_{q}}} \log \left|\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right|=\operatorname{tr}\left(G \partial_{\tau_{f_{q}}}\left(\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right)\right)=\operatorname{tr}\left(G \Lambda_{z}\right)  \tag{A.18}\\
& \partial_{\tau_{f_{r}}} \partial_{\tau_{f_{q}}} \log \left|\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right|\left.=\operatorname{tr}\left(\left(\partial_{\tau_{f_{r}}} G\right) \Lambda_{z^{\prime}}\right)+G\left(\partial_{\tau_{f_{r}}} \Lambda_{q}\right)\right)=\operatorname{tr}\left(-G \Lambda_{r} G \Lambda_{q}+G \Lambda_{r} \delta_{q r}\right)  \tag{A.19}\\
& \partial_{\eta_{u}} \log \left|\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right|=\operatorname{tr}\left(G \partial_{\eta_{u}}\left(\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right)\right)=\operatorname{tr}\left(G \widehat{e^{\eta(\boldsymbol{x}) \delta_{\boldsymbol{x} u}}}\right)  \tag{A.20}\\
& \partial_{\eta_{v}} \partial_{\eta_{u}} \log \left|\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right|=\operatorname{tr}\left(\left(\partial_{\eta_{v}} G\right) \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}} u}+G\left(\partial_{\eta_{v}} \widehat{e^{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}}\right)\right) \\
&=\operatorname{tr}\left(-G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} v}} G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}} u}+\delta_{v u} G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}}\right)  \tag{A.21}\\
&\left.\partial_{\eta_{u}} \partial_{\tau_{f_{q}} \log \left|\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right|}=\operatorname{tr}\left(\left(\partial_{\eta_{u}} G\right) \Lambda_{q}\right)\right)=\operatorname{tr}\left(-G e^{\widehat{(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G \Lambda_{q}\right)  \tag{A.22}\\
&=\partial_{\tau_{f_{q}}} \partial_{\eta_{u}} \log \left|\tilde{F}\left[\tau_{f}\right]+\widehat{e^{\eta(\boldsymbol{x})}}\right| \tag{A.23}
\end{align*}
$$

## A.3. The Determinant $\left|\Gamma_{\xi}\right|$

The determinant for the curvature of $\gamma_{\xi}$ is given by

$$
\begin{align*}
\left|\Gamma_{\xi}\left(\beta, \tau_{\beta}\right)\right|= & \left|\frac{\delta^{2} \gamma}{\delta \beta^{\dagger} \delta \beta}\right|\left|\frac{\delta^{2} \gamma}{\delta \tau_{\beta}^{\dagger} \delta \tau_{\beta}}-\frac{\delta^{2} \gamma}{\delta \tau_{\beta}^{\dagger} \delta \beta}\left(\frac{\delta^{2} \gamma}{\delta \beta^{\dagger} \delta \beta}\right)^{-1} \frac{\delta^{2} \gamma}{\delta \beta^{\dagger} \delta \tau_{\beta}}\right| \\
= & \operatorname{det}_{u v}\left[\widehat{\rho e^{\beta(\boldsymbol{z})}}{ }_{u v}+\left(\mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \mathcal{F}}\right)_{u v}\right]_{u v}^{\operatorname{det}_{u v}}\left[\frac{1}{2} \delta_{u v} \beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{\tau_{\beta} \delta_{x}} \mathcal{F} \beta+\left(\frac{1}{\sigma_{\beta}} \Delta^{\dagger} \Delta\right)_{u v}}\right. \\
& \left.-\left(\beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{\tau_{\beta} \delta_{u}} \mathcal{F}}\right)_{w}\left(\widehat{\rho e^{\beta(\boldsymbol{z})}}+\left(\mathcal{F}^{\dagger} \widehat{e^{-\tau_{\beta}} \mathcal{F}}\right)\right)_{w z}^{-1}\left(\beta^{\dagger} \mathcal{F}^{\dagger} \widehat{e^{\tau_{\beta} \delta_{v}} \mathcal{F}}\right)_{z}\right] \tag{A.24}
\end{align*}
$$

## A.4. The Determinant $\left|\Gamma_{\zeta}\right|$

The full expression for the determinant of the curvature of $\gamma_{\zeta}$ is:

$$
\begin{align*}
& \left\lvert\, \Gamma_{\zeta}\left(\eta, \tau_{f} \mid\right)=\left(\left|\frac{\delta^{2} \gamma}{\delta \tau_{f}^{\dagger} \delta \tau_{f}}\right|\left|\frac{\delta^{2} \gamma}{\delta \eta^{\dagger} \delta \eta}-\frac{\delta^{2} \gamma}{\delta \eta^{\dagger} \delta \tau_{f}}\left(\frac{\delta^{2} \gamma}{\delta \tau_{f}^{\dagger} \delta \tau_{f}}\right)^{-1} \frac{\delta^{2} \gamma}{\delta \tau_{f}^{\dagger} \delta \eta}\right|\right)\right. \\
& =\operatorname{det}_{u v}\left[\frac{1}{2} \operatorname{tr}\left(-G \Lambda_{x} G \Lambda_{y}+G \Lambda_{x} \delta_{u v}\right)+\frac{1}{2} \boldsymbol{y}^{\dagger}\left(2 G \Lambda_{u} G \Lambda_{v} G-G \Lambda_{u} \delta_{u v} G\right) \boldsymbol{y}+\left(\frac{1}{\sigma_{f}} \Delta^{\dagger} \Delta\right)_{u v}\right] \\
& \times \operatorname{det}_{u v}\left[\frac{1}{2} \operatorname{tr}\left(-G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} v}}+\delta_{u v} G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} v}}\right)\right. \\
& +\frac{1}{2} \boldsymbol{y}^{\dagger}\left(2 G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} v}} G-G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x}}} \delta_{u v} G\right) \boldsymbol{y}+\left(\frac{1}{\sigma_{\eta}} \nabla^{\dagger} \nabla\right)_{u v} \\
& -\left(-\frac{1}{2} \operatorname{tr}\left(G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G \Lambda_{w}\right)+\boldsymbol{y}^{\dagger} G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} u}} G \Lambda_{w} G \boldsymbol{y}\right) \times \\
& \times\left(\frac{1}{2} \operatorname{tr}\left(-G \Lambda_{w} G \Lambda_{z}+G \Lambda_{w} \delta_{w z}\right)+\frac{1}{2} \boldsymbol{y}^{\dagger}\left(2 G \Lambda_{w} G \Lambda_{z} G-G \Lambda_{w} \delta_{w z}\right) \boldsymbol{y}+\left(\frac{1}{\sigma_{f}} \Delta^{\dagger} \Delta\right)_{w z}\right)^{-1} \\
& \left.\times\left(-\frac{1}{2} \operatorname{tr}\left(G e^{\widehat{\eta(\boldsymbol{x})} \delta_{\boldsymbol{x} v}} G \Lambda_{z}\right)+\boldsymbol{y}^{\dagger} G e^{\widehat{\eta(\boldsymbol{x}) \delta_{\boldsymbol{x} v}}} G \Lambda_{z} G \boldsymbol{y}\right)\right] \tag{A.25}
\end{align*}
$$

## B. Details on the Computational Implementation

The implementation of algorithm 2 was done in Python ${ }^{1}$ Using the notation of 3, we modelled $\boldsymbol{x}, \boldsymbol{y}$, as 1 d -Numpy arrays with a length of $n_{\text {bins }}$.

The operators $B, F$ are represented as 2 d-Numpy arrays with a shape of ( $n_{\text {bins }}, n_{\text {bins }}$ ), i.e. matrices in $\mathbb{R}^{n_{\text {bins }} \times n_{\text {bins }}}$. Using a Fourier transform (i.e. a Discrete Fourier transform) to represent the $B$ and $F$ would lead to the question how to handle the resulting imaginary parts. Instead we choose to use the Hartley transform ( $(\operatorname{Har} 42)$ which is defined as

$$
\begin{equation*}
\mathcal{H}[f](q)=\frac{1}{2 \pi} \int \mathrm{~d} x(\sin (x q)+\cos (x q)) f(x) \tag{B.1}
\end{equation*}
$$

Therefore, $\mathcal{H}[f]=\mathfrak{R e}(\mathcal{F}[f])-\mathfrak{I m}(\mathcal{F}[f])$, and further $\mathcal{H}^{-1}=\mathcal{H}$. We thus implement the Hartley transform by calculating the discrete Hartley transform matrix, $\mathbb{H}=\mathfrak{R e}(\mathbb{F})-\mathfrak{I m}(\mathbb{F})$. Here, $\mathbb{F}$ denotes the discrete Fourier transform (DFT) matrix, for which we use an implementation provided by the SciPy Module. We can now give numerical representations for $B$ and $F$ via $\mathbb{H} \widehat{P_{\beta}} \mathbb{H}$ and $\mathbb{H} \widehat{P_{f}} \mathbb{H}$.

[^5]
## C. Scatter Plots for the Benchmark Data

Figure C.1.: Synthetic Datasets sampled via alg. 1 with $\varsigma^{2}=0.2$. Blue scatter plots indicate a true causal direction of $X \rightarrow Y$, red scatter plots the direction $Y \rightarrow X$.


Figure C.2.: Synthetic Datasets sampled via alg. 1 with $\varsigma^{2}=1.0$. Blue scatter plots indicate a true causal direction of $X \rightarrow Y$, red scatter plots the direction $Y \rightarrow X$.


Figure C.3.: Synthetic Datasets sampled via alg. 1 with $n_{\text {bins }}=16$. Blue scatter plots indicate a true causal direction of $X \rightarrow Y$, red scatter plots the direction $Y \rightarrow X$.


Figure C.4.: Synthetic Datasets sampled via alg. 1 with $n_{\text {bins }}=8$. Blue scatter plots indicate a true causal direction of $X \rightarrow Y$, red scatter plots the direction $Y \rightarrow X$.



Figure C.5.: Synthetic Datasets sampled via alg. 1 with $\tilde{N}=30$. Blue scatter plots indicate a true causal direction of $X \rightarrow Y$, red scatter plots the direction $Y \rightarrow X$.

| $\therefore$ |
| :---: |
| $\because$ |
| $\because$ |
| $\because$ |



| $\because \because$ |  |  | $\vdots$ |
| :--- | :--- | :--- | :--- |
| $\because$ | $\ddots$ | $\ddots$ |  |
|  | $\ddots$ | $\ddots$ |  |



Figure C.6.: Synthetic Datasets sampled via alg. 1 with $\tilde{N}=10$. Blue scatter plots indicate a true causal direction of $X \rightarrow Y$, red scatter plots the direction $Y \rightarrow X$.


## D. Description of the TCEP Benchmark Set

The following is a description for the TCEP benchmark set. We omit the 6 datasets which contain more than 2 variables as they have not been used in the benchmark. Scatter plots are shown in fig. D.1.

| dataset name | $X$ | $Y$ | source | causal direction |
| :---: | :---: | :---: | :---: | :---: |
| pair0001 | Altitude | Temperature | DWD | $X \rightarrow Y$ |
| pair0002 | Altitude | Precipitation | DWD | $X \rightarrow Y$ |
| pair0003 | Longitude | Temperature | DWD | $X \rightarrow Y$ |
| pair0004 | Altitude | Sunshine hours | DWD | $X \rightarrow Y$ |
| pair0005 | Age | Length | Abalone | $X \rightarrow Y$ |
| pair0006 | Age | Shell weight | Abalone | $X \rightarrow Y$ |
| pair0007 | Age | Diameter | Abalone | $X \rightarrow Y$ |
| pair0008 | Age | Height | Abalone | $X \rightarrow Y$ |
| pair0009 | Age | Whole weight | Abalone | $X \rightarrow Y$ |
| pair0010 | Age | Shucked weight | Abalone | $X \rightarrow Y$ |
| pair0011 | Age | Viscera weight | Abalone | $X \rightarrow Y$ |
| pair0012 | Age | Wage per hour | census income | $X \rightarrow Y$ |
| pair0013 | Displacement | Fuel consumption | auto-mpg | $X \rightarrow Y$ |
| pair0014 | Horse power | Fuel consumption | auto-mpg | $X \rightarrow Y$ |
| pair0015 | Weight | Fuel consumption | auto-mpg | $X \rightarrow Y$ |
| pair0016 | Horsepower | Acceleration | auto-mpg | $X \rightarrow Y$ |
| pair0017 | Age | Dividends from stocks | census income | $X \rightarrow Y$ |
| pair0018 | Age | Concentration GAG | GAGurine (from R package MASS) | $X \rightarrow Y$ |
| pair0019 | Current duration | Next interval | geyser | $X \rightarrow Y$ |
| pair0020 | Latitude | Temperature | DWD | $X \rightarrow Y$ |
| pair0021 | Longitude | Precipitation | DWD | $X \rightarrow Y$ |
| pair0022 | Age | Height | arrhythmia | $X \rightarrow Y$ |
| pair0023 | Age | Weight | arrhythmia | $X \rightarrow Y$ |
| pair0024 | Age | Heart rate | arrhythmia | $X \rightarrow Y$ |
| pair0025 | Cement | Compressive strength | concrete_data | $X \rightarrow Y$ |


| pair0026 | Blast furnace slag | Compressive strength | concrete_data | $X \rightarrow Y$ |
| :---: | :---: | :---: | :---: | :---: |
| pair0027 | Fly ash | Compressive strength | concrete_data | $X \rightarrow Y$ |
| pair0028 | Water | Compressive strength | concrete_data | $X \rightarrow Y$ |
| pair0029 | Superplasticizer | Compressive strength | concrete_data | $X \rightarrow Y$ |
| pair0030 | Coarse aggregate | Compressive strength | concrete_data | $X \rightarrow Y$ |
| pair0031 | Fine aggregate | Compressive strength | concrete_data | $X \rightarrow Y$ |
| pair0032 | Age | Compressive strength | concrete_data | $X \rightarrow Y$ |
| pair0033 | Alcohol consumption | Mean corpuscular volume | liver disorders | $X \rightarrow Y$ |
| pair0034 | Alcohol consumption | Alkaline phosphotase | liver disorders | $X \rightarrow Y$ |
| pair0035 | $\begin{aligned} & \text { Alcohol con- } \\ & \text { sumption } \end{aligned}$ | Alanine aminotransferase | liver disorders | $X \rightarrow Y$ |
| pair0036 | Alcohol consumption | Aspartate aminotransferase | liver disorders | $X \rightarrow Y$ |
| pair0037 | Alcohol consumption | Gammaglutamyl transpeptdase | liver disorders | $X \rightarrow Y$ |
| pair0038 | Age | Body mass index | pima indian diabetes | $X \rightarrow Y$ |
| pair0039 | Age | Serum insulin | pima indian diabetes | $X \rightarrow Y$ |
| pair0040 | Age | Diastolic blood pressure | pima indian diabetes | $X \rightarrow Y$ |
| pair0041 | Age | Plasma glucose concentration | pima indian diabetes | $X \rightarrow Y$ |
| pair0042 | Day of the year | Temperature | B.Janzing | $X \rightarrow Y$ |
| pair0043 | Temperature at t | Temperature at t+1 | ncep-ncar | $X \rightarrow Y$ |
| pair0044 | Pressure at t | Pressure at t+1 | ncep-ncar | $X \rightarrow Y$ |
| pair0045 | Sea level pressure at t | Sea level pressure at $\mathrm{t}+1$ | ncep-ncar | $X \rightarrow Y$ |
| pair0046 | Relative humidity at t | Relative humidity at $\mathrm{t}+1$ | ncep-ncar | $X \rightarrow Y$ |
| pair0047 | Number of cars | Type of day | traffic | $Y \rightarrow X$ |
| pair0048 | Indoor temperature | Outdoor temperature | Hipel \& Mcleod | $Y \rightarrow X$ |
| pair0049 | Ozone concentration | Temperature | Bafu | $Y \rightarrow X$ |


| pair0050 | Ozone concentration | Temperature | Bafu | $Y \rightarrow X$ |
| :---: | :---: | :---: | :---: | :---: |
| pair0051 | Ozone concentration | Temperature | Bafu | $Y \rightarrow X$ |
| pair0056 | Female life expectancy, 20002005 | Latitude | UNdata | $Y \rightarrow X$ |
| pair0057 | Female life expectancy, 19952000 | Latitude | UNdata | $Y \rightarrow X$ |
| pair0058 | Female life expectancy, 19901995 | Latitude | UNdata | $Y \rightarrow X$ |
| pair0059 | Female life expectancy, 19851990 | Latitude | UNdata | $Y \rightarrow X$ |
| pair0060 | Male life expectancy, 2000-2005 | Latitude | UNdata | $Y \rightarrow X$ |
| pair0061 | Male life expectancy, 1995-2000 | Latitude | UNdata | $Y \rightarrow X$ |
| pair0062 | Male life expectancy, 1990-1995 | Latitude | UNdata | $Y \rightarrow X$ |
| pair0063 | Male life expectancy, 1985-1990 | Latitude | UNdata | $Y \rightarrow X$ |
| pair0064 | Drinking water access | Infant mortality | UNdata | $X \rightarrow Y$ |
| pair0065 | Stock return of Hang Seng Bank | Stock return of HSBC Hldgs | Yahoo database | $X \rightarrow Y$ |
| pair0066 | Stock return of Hutchison | Stock return of Cheung kong | Yahoo database | $X \rightarrow Y$ |
| pair0067 | Stock return of Cheung kong | Stock return of Sun Hung Kai Prop. | Yahoo database | $X \rightarrow Y$ |
| pair0068 | Bytes sent | Open http connections | P. Stark \& Janzing | $Y \rightarrow X$ |
| pair0069 | Inside temperature | Outside temperature | J.M. Mooij | $Y \rightarrow X$ |
| pair0070 | Parameter | Answer | Armann $\quad \& ~$ Buelthoff | $X \rightarrow Y$ |
| pair0072 | Sunspots | Global mean temperature | sunspot data | $X \rightarrow Y$ |
| pair0073 | CO2 emissions | Energy use | UNdata | $Y \rightarrow X$ |
| pair0074 | GNI per capita | Life expectancy | UNdata | $X \rightarrow Y$ |


| pair0075 | Under-5 mortality rate | GNI per capita | UNdata | $Y \rightarrow X$ |
| :---: | :---: | :---: | :---: | :---: |
| pair0076 | Population growth | Food consumption growth | Food and Agriculture Organization of the Unite... | $X \rightarrow Y$ |
| pair0077 | Temperature | Solar radiation | B. Janzing | $Y \rightarrow X$ |
| pair0078 | PPFD | Net Ecosystem Productivity | Moffat A.M. | $X \rightarrow Y$ |
| pair0079 | Net Ecosystem Productivity | Diffuse PPFDdif | Moffat A.M. | $Y \rightarrow X$ |
| pair0080 | Net Ecosystem Productivity | Direct PPFDdir | Moffat A.M. | $Y \rightarrow X$ |
| pair0081 | Temperature | Local CO2 flux, BE-Bra | Mahecha, M. | $X \rightarrow Y$ |
| pair0082 | Temperature | Local CO2 flux, DE-Har | Mahecha, M. | $X \rightarrow Y$ |
| pair0083 | Temperature | Local CO2 flux, US-PFa | Mahecha, M. | $X \rightarrow Y$ |
| pair0084 | Employment | Population | spatialeconometrics.com | $Y \rightarrow X$ |
| pair0085 | Time of measurement | Protein content of milk | maths.lancs.ac.uk | $X \rightarrow Y$ |
| pair0086 | Size of apartment | Monthly rent | J.M. Mooij | $X \rightarrow Y$ |
| pair0087 | Temperature | Total snow | Snowfall in Whistler, from www.mldata.org | $X \rightarrow Y$ |
| pair0088 | Age | Relative spinal bone mineral density | bone dataset of R ElemStatLearn package | $X \rightarrow Y$ |
| pair0089 | root decomposition Oct (grassl) | root decomposition Oct (grassl) | Solly et al (2014). Plant and Soil, 382(12), ... | $Y \rightarrow X$ |
| pair0090 | root decomposition Oct (forest) | root decomposition Oct (forest) | Solly et al (2014). Plant and Soil, 382(12), ... | $Y \rightarrow X$ |
| pair0091 | clay cont. in soil (forest) | soil moisture | Solly et al (2014). Plant and Soil, 382(12), ... | $X \rightarrow Y$ |


| pair0092 | organic carbon in soil (forest) | clay cont. in soil (forest) | Solly et al (2014). Plant and Soil, 382(12), ... | $Y \rightarrow X$ |
| :---: | :---: | :---: | :---: | :---: |
| pair0093 | precipitation | runoff | MOPEX | $X \rightarrow Y$ |
| pair0094 | hour of day | temperature | S. Armagan Tarim | $X \rightarrow Y$ |
| pair0095 | hour of day | electricity load | S. Armagan Tarim | $X \rightarrow Y$ |
| pair0096 | temperature | electricity load | S. Armagan Tarim | $X \rightarrow Y$ |
| pair0097 | speed at the beginning | speed at the end | D. Janzing | $X \rightarrow Y$ |
| pair0098 | speed at the beginning | speed at the end | D. Janzing | $X \rightarrow Y$ |
| pair0099 | language test score | social-economic status family | nlschools <br> dataset of R | $Y \rightarrow X$ |
| pair0100 | cycle time of CPU | performance | MASS package cpus dataset of R MASS package | $X \rightarrow Y$ |
| pair0101 | grey value of a pixel | brightness of the screen | D. Janzing | $X \rightarrow Y$ |
| pair0102 | position of a ball | time for passing a track segment | D. Janzing | $X \rightarrow Y$ |
| pair0103 | position of a ball | time for passing a track segment | D. Janzing | $X \rightarrow Y$ |
| pair0104 | time for passing <br> 1. segment | time for passing <br> 2. segment | D. Janzing | $X \rightarrow Y$ |
| pair0106 | time required for one round | voltage | D. Janzing | $Y \rightarrow X$ |
| pair0107 | strength of contrast | answer correct or not | Schuett, edited by D. Janzing | $X \rightarrow Y$ |
| pair0108 | time for $1 / 6$ rotation | temperature | D. Janzing | $Y \rightarrow X$ |

Figure D.1.: TCEP benchmark dataset. Blue scatter plots indicate a true causal direction of $X \rightarrow Y$, red scatter plots the direction $Y \rightarrow X$.


## Statement of authorship

I hereby declare that I completed this thesis on my own and that information which has been directly or indirectly taken from other sources has been noted as such.

Munich, September 12, 2018


[^0]:    ${ }^{1}$ Using the standard notation for directed graphs here, where an edge $(a, b)$ indicates the direction " $a$ to $b$ ".

[^1]:    ${ }^{1}$ This property is called stationarity in the context of stochastic processes (see e.g.

    ```
    Cha16
    ```

[^2]:    ${ }^{2}$ This type of matrix, i.e. the evaluation of covariance or kernel at certain positions, is sometimes called a Gram matrix.

[^3]:    * As we draw the number of samples from Poisson distribution in each bin, we cannot deterministically control the total number of samples

[^4]:    $2^{2}$ https://github.com/MKurthen/BayesianCausalInference

[^5]:    ${ }^{1}$ provided at https://github.com/MKurthen/BayesianCausalInference

