
Design of Dynamic Experiments for Black-Box Model Discrimination

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Abstract Diverse domains of science and engineering require and use mechanistic mathematical models, e.g. systems of differential algebraic equations. Such models often contain uncertain parameters to be estimated from data. Consider a dynamic model discrimination setting where we wish to choose: (i) what is the best mechanistic, time-varying model and (ii) what are the best model parameter estimates. These tasks are often termed model discrimination/selection/validation/verification. Typically, several rival mechanistic models can explain data, so we incorporate available data and also run new experiments to gather more data. Design of dynamic experiments for model discrimination helps optimally collect data. For rival mechanistic models where we have access to gradient information, we extend existing methods to incorporate a wider range of problem uncertainty and show that our proposed approach is equivalent to historical approaches when limiting the types of considered uncertainty. We also consider rival mechanistic models as dynamic *black boxes* that we can evaluate, e.g. by running legacy code, but where gradient or other advanced information is unavailable. We replace these black-box models with Gaussian process surrogate models and thereby extend the model discrimination setting to additionally incorporate rival black-box model. We also explore the consequences of using Gaussian process surrogates to approximate gradient-based methods.

Keywords Experimental design · Model discrimination · Gaussian processes

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1 Introduction

Mathematical models based on mechanistic assumptions are common in both engineering, e.g. pharmaceutical manufacturing (Rathore and Winkle, 2009; Chen et al., 2019), and in the natural sciences, e.g. climate forecasting and high-energy physics (Altarelli, 2014; Mashnik, 2010; Flato et al., 2014). Researchers and engineers often hypothesise several such mechanistic models to explain the underlying system behaviour. Model discrimination differentiates between these rival hypotheses, i.e. discards inaccurate models (Hunter and Reiner, 1965; Box and Hill, 1967). Often the available experimental data do not allow for discrimination between the rival models, i.e. all models sufficiently explain the available data, at least within data accuracy, e.g. Joy et al. (2019). This situation requires more data, e.g. by running additional experiments. But experiments are expensive, so we wish to discriminate between the models with as few additional experiments as possible (Walz et al., 2018, 2020).

As an example of requiring effective model discrimination, consider the parametric, mechanistic models used in the heavily regulated pharmaceutical manufacturing industry. The US Food & Drug Administration (FDA) and European Medicines Agency (EMA) regulate the development, manufacturing and marketing of drugs and other pharmaceutical products. Drugs were not always heavily regulated: modern-day regulations result from historical failures to protect patients (Wax, 1995). Historical drug scandals led to improved regulations and thorough testing of new pharmaceutical products.

To predict and prevent such tragedies while still encouraging innovation, both the FDA and EMA incorporate *quality-by-design* (QbD) into their regulations (Rathore and Winkle, 2009). QbD, which guarantees product quality through process design rather than by simply testing the final product, requires a deep understanding of a product and its manufacturing process. Typically, these QbD guidelines for pharmaceutical manufacturing require mechanistic understanding of the final product quality (U.S. Food & Drug Administration, 2009). These mechanistic models are typically parametric, with model parameters tuned to fit the experimental data, e.g. as in Lakerveld et al. (2013).

To illustrate model discrimination, consider three rival, linear continuous-time models \mathcal{M}_1 , \mathcal{M}_2 and \mathcal{M}_3 (Bania, 2019)

$$\mathcal{M}_i : \begin{cases} \left. \frac{dx}{dt} \right|_t = \mathbf{A}_i x(t) + \mathbf{B}_i u(t) + \mathbf{C}_i w(t), \\ y_k = [1, 0, \dots, 0] x(t_k) + 0.05 \cdot v_k, \end{cases} \quad (1)$$

where $x(t)$ denotes the state at time t , y_k denotes measurements of a subset of the states at discrete time points t_k , $k = 1, \dots, T$, $u(t)$ denotes the control input, $w(t)$ is the process noise and $v(t)$ is the measurement noise. Fig. 1 shows the model predictions for three different piece-wise constant control signals $u(t) \in [-1, 1]$, with matrices \mathbf{A}_i , \mathbf{B}_i and \mathbf{C}_i defined in Appendix A. The initial state is $\mathbf{x}_0^{(i)} = [0, \dots, 0]^\top$, and the process noise $w(t) \sim \mathcal{N}(0, 1)$ and the measurement noise $v_k \sim \mathcal{N}(0, 1)$ are Gaussian distributed.

Fig. 1a shows a step input, and Fig. 1c uniformly distributed random inputs. Fig. 1b and Fig. 1d show the corresponding approximate marginal predictive distributions over time. The different models \mathcal{M}_1 , \mathcal{M}_2 and \mathcal{M}_3 yield different inter-

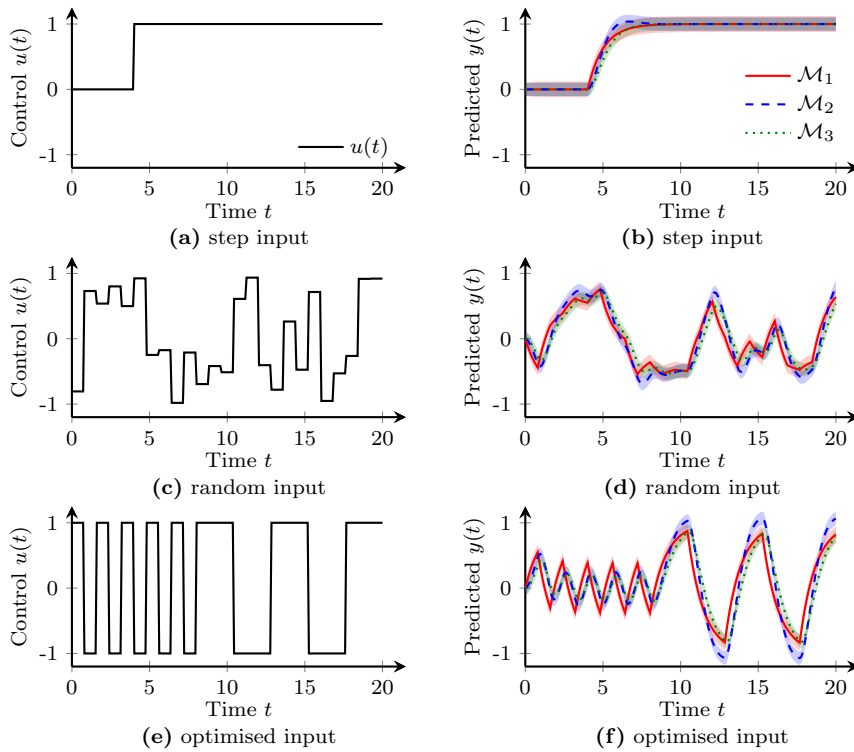


Fig. 1 The three Bania (2019) models \mathcal{M}_1 , \mathcal{M}_2 and \mathcal{M}_3 , defined in (1). (a), (c) and (e) show three different piece-wise constant control input as function of time t , and (b), (d) and (e) show the corresponding model predictions, with mean and two standard deviations.

pretations of underlying system mechanisms. But it is non-trivial to find control inputs yielding sufficiently different predictions to allow for model discrimination.

Observe in Fig. 1b and Fig. 1d that naively applying a step or random control input may not produce data resolving the model discrimination problem. A naive control input may also result in violations of system safety constraints. Fig. 1e shows an example of a different control signal (developed using the methods presented in this paper). This control signal is optimised to yield sufficiently large model predictions differences (shown in Fig. 1f) for model discrimination. More than one experiment may be required to discriminate between models, and a systematic approach is necessary to minimise the number of experiments.

The example provided by Bania (2019) gives easy access to model gradients. But many industrially relevant models cannot be written in closed form, e.g. if they are implemented in legacy code. From an optimisation perspective, these models are effectively black boxes. We can evaluate the models, but the gradients are not readily available. Specialist software can generate the gradients automatically, but this poses substantial requirements on the implementation of the specialist models (Naumann, 2012). For model discrimination, we wish to be agnostic to the software implementation or model type, since this flexibly (i) allows for faster model prototyping and development, and (ii) satisfies the personal preferences of researchers and engineers. This paper finds optimal control inputs and other

experimental conditions, e.g. the initial state $x(0)$, to discriminate dynamic white- or black-box models under uncertainty and subject to constraints.

2 Design of Dynamic Experiments for Model Discrimination

We consider state-space formulations of dynamic models (Hangos *et al.*, 2004, Ch. 3). Most of this paper focuses on discrete-time models, but Appendix C additionally shows how to extend our contributions to continuous-time models.

Let $\mathcal{T} = \{t_0, \dots, t_T\}$ denote an ordered set of discrete time instances. These time instances are indexed by $k = 0, \dots, T$, with $k = 0$ the starting time step of an experiment and $k = T$ the final time step. For simplicity, we write $k \in \mathcal{T}$. Assume the system at time step k is in a state $\mathbf{x}_k = x(t_k) \in \mathbb{R}^{D_x}$. The rate of change of the state is a function of the current state \mathbf{x}_k , a control input $\mathbf{u}_k = u(t_k) \in \mathbb{R}^{D_u}$, and some process noise \mathbf{w}_k . The models describing the rate of change are parameterised by $\boldsymbol{\theta} \in \mathbb{R}^{D_\theta}$. The control signal $u(t)$ is commonly discretised and piece-wise constant in order to make optimisation feasible in dynamic problems with continuous-time models, e.g. in Schultz *et al.* (2019). In general, we cannot observe all dimensions of the state \mathbf{x}_k . This is commonly the case for pharmaceutical manufacturing where the models predict different chemical concentrations, only some of which can be measured. Let $\mathbf{z}_k = z(t_k) \in \mathbb{R}^{D_z}$ denote the observed state. We take measurements $\mathbf{y}_k \in \mathbb{R}^{D_z}$ of the observed states \mathbf{z}_k at discrete time instances $k \in \mathcal{T}_{\text{meas}}$. Measurement noise \mathbf{v}_k corrupts the measurements.

Appendix C and Appendix D, respectively, extend our work to continuous-time models and discrete-time models with a Δ -transition, which follows from an Euler discretisation of continuous-time dynamics (Atkinson *et al.*, 2009, Ch. 2).

Let $\mathbf{x}_{0:T}$, $\mathbf{z}_{1:T}$ and $\mathbf{u}_{0:T-1}$ denote the sequences of states, observed states and control inputs, respectively. The inputs to the system have dimensionality $T \times D_u + D_\theta$, and the outputs have dimensionality $T \times D_z$. Assume M rival state space models $\mathcal{M}_1, \dots, \mathcal{M}_M$. Model \mathcal{M}_i specifies a state transition, observation and measurement model given a state $\mathbf{x}_k^{(i)}$, control input \mathbf{u}_k at time step k , and model parameter vector $\boldsymbol{\theta}_i$

$$\mathcal{M}_i : \begin{cases} \mathbf{x}_k^{(i)} = f_i(\mathbf{x}_{k-1}^{(i)}, \mathbf{u}_{k-1}, \boldsymbol{\theta}_i) + \mathbf{w}_{k-1}^{(i)}, & \text{(State transition)} \\ \mathbf{z}_k^{(i)} = \mathbf{H}_i \mathbf{x}_k^{(i)}, & \text{(Observed states)} \\ \mathbf{y}_k^{(i)} = \mathbf{z}_k^{(i)} + \mathbf{v}_k^{(i)}, & \text{(Noisy measurement)} \end{cases} \quad (2)$$

where f_i is the transition function, and \mathbf{H}_i is a matrix selecting the observed states. The number of states $D_{x,i}$ may differ between models \mathcal{M}_i , but the number of observed states $D_{z,i} = D_z$ for all models.

We assume the observed state \mathbf{z}_k is a subset (or linear combination) of the state \mathbf{x}_k , i.e., $z(t) = \mathbf{H}x(t)$. Since the transition function f can be any nonlinear function, this is a minor restriction to the types of models we consider for domains such as pharmaceuticals and chemical manufacturing. Moreover, the methods described in this manuscript could be extended to the more general case of a nonlinear relationship $z(t) = g(x(t))$, where g is known.

The fundamental principle of sequential experimental design for model discrimination says to select the next experimental point where the model predictions

| Ref. | (a) | (b) | (c) | (d) | (e) | (f) |
|--------------------------------------|-----|-----|-----|-----|-----|-----|
| Nonlinear f | ✓ | ✓ | | ✓ | (✓) | ✓ |
| Discrete-time models | | | ✓ | | ✓ | ✓ |
| Continuous-time models | ✓ | ✓ | | ✓ | | ✓ |
| Black-box models | | | | ✓ | | ✓ |
| Measurement noise | ✓ | ✓ | (✓) | ✓ | ✓ | ✓ |
| Process noise | | | ✓ | | ✓ | ✓ |
| Uncertain \mathbf{x}_0 | | | ✓ | ✓ | | ✓ |
| Uncertain \mathbf{u}_k | | | | | | ✓ |
| Uncertain $\boldsymbol{\theta}$ | ✓ | (✓) | | ✓ | | ✓ |
| Optimise \mathbf{x}_0 | ✓ | ✓ | | | | ✓ |
| Optimise $\mathbf{u}_{0:T-1}$ | ✓ | ✓ | ✓ | ✓ | ✓ | ✓ |
| Optimise $\mathcal{T}_{\text{meas}}$ | ✓ | ✓ | | | | ✓ |
| Path constraints | | (✓) | (✓) | | | ✓ |

Table 1 References: (a) Chen and Asprey (2003), (b) Skanda and Lebedez (2013), (c) Cheong and Manchester (2014), (d) Streif et al. (2014), (e) Bania (2019), (f) This work. Bracketed check marks: Bania (2019) mention how their approach can be extended to non-linear transition functions f ; Cheong and Manchester (2014) has one noise signal that affects both states and measurements; Skanda and Lebedez (2013) use a robust problem formulation instead of marginalising out the model parameters; Skanda and Lebedez (2013) and Cheong and Manchester (2014) solve the design of experiments optimisation problem subject to path constraints that do not account for predicted state uncertainty.

differ the most (Hunter and Reiner, 1965). Due to the multiple sources of uncertainty, e.g. measurement noise, this translates into maximising the divergence between the models’ marginal predictive distributions. The problems are (i) approximating the models’ marginal predictive distributions, and (ii) constructing and maximising a predictive distribution divergence measure. There are two main approaches to solve these problems: the analytical approach, and the data-driven approach.

The most common analytical approach assumes all uncertainty sources are Gaussian distributed and uses a first-order Taylor expansion around the input mean to propagate the Gaussian input distribution $p(\mathbf{x}_{k-1}, \mathbf{u}_{k-1}, \boldsymbol{\theta})$ to a Gaussian output distribution $p(\mathbf{x}_k | \dots)$ (Box and Hill, 1967; MacKay, 1992; Chen and Asprey, 2003). Appendix B offers a more detailed overview. Several closed-form divergence measures exist for Gaussian marginal predictive distributions (Box and Hill, 1967; Buzzi-Ferraris et al., 1990; Michalik et al., 2010; Olofsson et al., 2019). The analytical approach may be computationally efficient, but limited in the sense that (i) it uses linear and Gaussian approximations to handle uncertainty, and (ii) it requires derivative information.

The data-driven approach uses Monte Carlo techniques to approximate the marginal predictive distributions’ divergence (Ryan et al., 2016). Variational techniques can also be used (Foster et al., 2019). This approach is flexible, because it is agnostic to the model structure (and indifferent to the presence or absence of function gradients). It may also be more accurate given enough samples, since it does not rely on linearising the models or Gaussian distributed uncertainty. However, it is computationally expensive, does not scale well with the problem size and makes optimisation difficult. As mentioned, the dimensionality of the input and output spaces for the rival models can grow large when we discretise in time. The

computational cost associated with solving design of dynamic experiments problems using Monte Carlo techniques would likely become insurmountable. Streif *et al.* (2014) and Paulson *et al.* (2019) use polynomial chaos methods for design of dynamic experiments, but must still rely on expensive approximations of the predictive distributions' divergence.

Olofsson *et al.* (2018) propose a hybrid approach for black-box model discrimination. They replace the black-box models with Gaussian process (GP) surrogates. Training data for the GP surrogates are collected by evaluating the original black-box model at sampled locations. The GP surrogates are then used in an analytical fashion. This approach is computationally relatively cheap and can accommodate black-box models. However, naively applying the Olofsson *et al.* (2018) approach to design of dynamic experiments is intractable, since it would operate directly on the mapping from $(\mathbf{x}_0, \mathbf{u}_{0:T-1}, \mathcal{T}_{\text{meas}})$ to $\mathbf{z}_{1:T}$. The GP surrogate input dimensionality then scales as $\mathcal{O}(D_x + D_u \times T + D_\theta)$ and output dimensionality as $\mathcal{O}(D_z \times T)$. Even for small T , e.g. $T = 20$, the dimensionality of the input space is too high for accurate GP inference. The number of GP surrogates would have to equal the output dimensionality, which would be very expensive memory-wise.

Columns (a)–(e) of Table 1 summarise several proposed methods in design of dynamic experiments for model discrimination. This paper extends both the analytical approach of Chen and Asprey (2003) and the Olofsson *et al.* (2018) GP approach, to design *dynamic* experiments for discriminating black-box models. Column (f) shows the novelty of our proposed approach compared to existing literature, e.g. analytically accommodating black-box models as well as accounting for more possible uncertainty types. The proposed approach's GP surrogate input dimensionality scales as $\mathcal{O}(D_x + D_u + D_\theta)$ and output dimensionality as $\mathcal{O}(D_x)$.

2.1 Gaussian Process Regression

A Gaussian process (GP) is a collection of random variables, any finite subset of which is jointly Gaussian distributed. A GP is completely specified by a mean function m and covariance function k (Rasmussen and Williams, 2006). Assume observations $\mathbf{q} = [q_1, \dots, q_N]^\top$ at locations $\mathbf{R} = [\mathbf{r}_1, \dots, \mathbf{r}_N]^\top$, with $q_n \sim \mathcal{N}(g(\mathbf{r}_n), \sigma_\eta^2)$, of a latent function $g : \mathbb{R}^{D_r} \rightarrow \mathbb{R}$ and zero-mean i.i.d. Gaussian observation noise with variance σ_η^2 . Let ψ denote the GP hyperparameters, which consist of the observation noise variance σ_η^2 and covariance function k 's parameters. GP regression computes a predictive distribution $g(\mathbf{r}) | \mathbf{R}, \mathbf{q}, \psi \sim \mathcal{N}(\mu(\mathbf{r}_*), \sigma^2(\mathbf{r}_*))$ at an arbitrary test point \mathbf{r}_* , with

$$\mu(\mathbf{r}_*) = m(\mathbf{r}_*) + \mathbf{k}^\top (\mathbf{K} + \sigma_\eta^2 \mathbf{I})^{-1} (\mathbf{q} - \mathbf{m}), \quad (3a)$$

$$\sigma^2(\mathbf{r}_*) = k(\mathbf{r}_*, \mathbf{r}_*) - \mathbf{k}^\top (\mathbf{K} + \sigma_\eta^2 \mathbf{I})^{-1} \mathbf{k}, \quad (3b)$$

where $[\mathbf{K}]_{j,\ell} = k(\mathbf{r}_j, \mathbf{r}_\ell)$, $[\mathbf{k}]_j = k(\mathbf{r}_*, \mathbf{r}_j)$ and $[\mathbf{m}]_j = m(\mathbf{r}_j)$. The hyperparameters ψ are typically learnt by maximising the marginal likelihood $p(\mathbf{q} | \mathbf{R}, \psi)$.

For vector-valued functions, i.e. vector fields, $g : \mathbb{R}^{D_r} \rightarrow \mathbb{R}^{D_q}$, we can place independent GP priors $\mathcal{GP}(m_{(d)}, k_{(d)})$ on each target dimension $d = 1, \dots, D_q$, with corresponding covariance function hyperparameters $\psi_{(d)}$. This yields the posterior distribution $g(\mathbf{r}_*) | \mathbf{R}, \mathbf{Q}, \Psi \sim N(M(\mathbf{r}_*), \Sigma(\mathbf{r}_*))$, with

$$M(\mathbf{r}_*) = [\mu_{(1)}(\mathbf{r}_*), \dots, \mu_{(D_q)}(\mathbf{r}_*)]^\top, \quad (4a)$$

$$\Sigma(\mathbf{r}_*) = \text{diag}(\sigma_{(1)}^2(\mathbf{r}_*), \dots, \sigma_{(D_q)}^2(\mathbf{r}_*)), \quad (4b)$$

where $\mu_{(d)}(\mathbf{r}_*)$ and $\sigma_{(d)}^2(\mathbf{r}_*)$ are the mean and variance given by (3) of the d^{th} GP posterior, $\mathbf{Q} = [\mathbf{q}_{(1)}, \dots, \mathbf{q}_{(D_q)}]$, with $\mathbf{q}_{(d)}$ dimension d of all observations, and $\Psi = \{\psi_{(1)}, \dots, \psi_{(D_q)}\}$ the joint set of hyperparameters. Independent GP priors can accommodate different covariance functions and hyperparameters for different target dimensions.

If the input is uncertain, e.g. $\mathbf{r}_* \sim \mathcal{N}(\boldsymbol{\mu}_*, \boldsymbol{\Sigma}_*)$, the posterior distribution needs to account for the added uncertainty. The exact posterior distribution is generally intractable, but can be approximated, e.g. using moment matching

$$p(g(\mathbf{r}_*) | \mathbf{R}, \mathbf{Q}, \Psi, \boldsymbol{\mu}_*, \boldsymbol{\Sigma}_*) \approx \mathcal{N}(\mathbb{E}_{\mathbf{r}_*} [M(\mathbf{r}_*)], \mathbb{E}_{\mathbf{r}_*} [\Sigma(\mathbf{r}_*)] + \mathbb{V}_{\mathbf{r}_*} [M(\mathbf{r}_*)]). \quad (5)$$

Computing the marginal mean and variance in (5) is also often intractable, with the exception of some special cases (Quiñonero-Candela et al., 2003; Deisenroth et al., 2009). An alternative approximation uses a first-order Taylor expansion of the predictive mean and variance in (3) with respect to \mathbf{r} , and propagates the input distribution through the linearisation $p(g(\mathbf{r}_*) | \mathbf{R}, \mathbf{Q}, \Psi, \boldsymbol{\mu}_*, \boldsymbol{\Sigma}_*) \approx \mathcal{N}(M(\boldsymbol{\mu}_*), \nabla_{\mathbf{r}} M \boldsymbol{\Sigma}_* \nabla_{\mathbf{r}} M^\top + \Sigma(\boldsymbol{\mu}_*))$ where $\nabla_{\mathbf{r}} M = \partial M(\mathbf{r}) / \partial \mathbf{r} |_{\mathbf{r}=\boldsymbol{\mu}_*}$. The linearisation is analytically tractable. By construction, the covariance $\nabla_{\mathbf{r}} M \boldsymbol{\Sigma}_* \nabla_{\mathbf{r}} M^\top + \Sigma(\boldsymbol{\mu}_*)$ is positive definite.

3 Experimental Design for Model Discrimination Under Uncertainty

We consider the general model formulation in (2), which makes it possible to extend existing methods, particularly that of Chen and Asprey (2003). We create a common optimal experimental design framework for discrimination between rival dynamic models with (i) analytical or black-box transition functions f_i , (ii) continuous- or discrete-time models, and (iii) multiple types of uncertainty.

In (2), $\mathbf{w}_k^{(i)}$ is zero-mean Gaussian-distributed process noise and $\mathbf{v}_k^{(i)} \sim \mathcal{N}(\mathbf{0}, \boldsymbol{\Sigma}_y)$ is independent and identically distributed measurement noise. The process noise $\mathbf{w}_k^{(i)}$ has known covariance $\boldsymbol{\Sigma}_{x,i}$. We assume the measurement noise covariance $\boldsymbol{\Sigma}_y$ is known but may be different ($\boldsymbol{\Sigma}_y = \boldsymbol{\Sigma}_y^{(i)}$) for different models \mathcal{M}_i ¹. Apart from uncertainty due to process noise and measurement noise, we may have uncertainty in the control input, initial state and model parameters.

The control input $\mathbf{u}_k \sim \mathcal{N}(\hat{\mathbf{u}}_k, \boldsymbol{\Sigma}_{u,k})$ at time step k is Gaussian distributed with mean given by a user-specified desired control input $\hat{\mathbf{u}}_k$ and covariance $\boldsymbol{\Sigma}_{u,k}$. The control covariance $\boldsymbol{\Sigma}_{u,k} = \boldsymbol{\Sigma}_{u,k}^{(i)}$ may be model-dependent. For simplicity, let the control inputs \mathbf{u}_k be piece-wise constant and the control covariance constant $\boldsymbol{\Sigma}_{u,k} = \boldsymbol{\Sigma}_u$. Simple extensions of the framework could accommodate control inputs described e.g. by piece-wise polynomials or time-dependent control covariance. Let $\hat{\mathbf{u}}_{0:T-1} = \{\hat{\mathbf{u}}_0, \dots, \hat{\mathbf{u}}_{T-1}\}$ denote the sequence of user-specified control inputs.

The initial state $\mathbf{x}_0^{(i)} \sim \mathcal{N}(\boldsymbol{\mu}_0^{(i)}(\hat{\mathbf{x}}_0), \boldsymbol{\Sigma}_0^{(i)})$ depends on some user-specified initial state settings $\hat{\mathbf{x}}_0$ common for all models. The vector $\hat{\mathbf{x}}_0$ is the desired initial state. Model parameters $\boldsymbol{\theta}_i \sim \mathcal{N}(\hat{\boldsymbol{\theta}}_i, \boldsymbol{\Sigma}_{\theta,i})$ are Gaussian distributed with mean given

¹ Note that model-specific measurement noise covariances interferes with some design criterion definitions.

by the maximum *a posteriori* parameter estimate $\hat{\boldsymbol{\theta}}_i$. A Laplace approximation computes the model parameter covariance $\boldsymbol{\Sigma}_{\theta,i}$ (MacKay, 2003, Ch. 27).

The initial state settings $\hat{\mathbf{x}}_0$ and the control inputs $\hat{\mathbf{u}}_{0:T-1}$ determine the experimental outcomes. For continuous-time models \mathcal{M}_i , we may also want to optimise the measurement time points $\mathcal{T}_{\text{meas}}$. Hence, we formulate the optimisation problem of design of dynamic experiments for model discrimination as

$$\begin{aligned} & \arg \max_{\substack{\hat{\mathbf{x}}_0, \hat{\mathbf{u}}_{0:T-1} \\ \mathcal{T}_{\text{meas}}} } \sum_{k \in \mathcal{T}_{\text{meas}}} D(\mathbf{y}_k^{(1)}, \dots, \mathbf{y}_k^{(M)}) \\ & \text{s.t. } \forall k \in \{1, \dots, T\}, \forall i \in \{1, \dots, M\} : \\ & \mathcal{M}_i : \begin{cases} \mathbf{x}_k^{(i)} = f_i(\mathbf{x}_{k-1}^{(i)}, \mathbf{u}_{k-1}, \boldsymbol{\theta}_i) + \mathbf{w}_{k-1}^{(i)}, \\ \mathbf{z}_k^{(i)} = \mathbf{H}_i \mathbf{x}_k^{(i)}, \\ \mathbf{y}_k^{(i)} = \mathbf{z}_k^{(i)} + \mathbf{v}_k, \end{cases} \quad (6) \\ & C_{x_0}(\hat{\mathbf{x}}_0) \geq \mathbf{0}, \quad \left| \quad C_x(\mathbf{x}_k^{(i)}) \geq \mathbf{0}, \quad \left| \quad C_{\mathcal{T}}(\mathcal{T}_{\text{meas}}) \geq \mathbf{0}. \right. \right. \\ & C_u(\hat{\mathbf{u}}_k) \geq \mathbf{0}, \quad \left| \quad C_z(\mathbf{z}_k^{(i)}) \geq \mathbf{0}, \quad \left| \right. \right. \end{aligned}$$

$D(\cdot)$ is the design criterion, i.e. a divergence measure between the predictive distributions. C_{x_0} , C_u , C_x , C_z and $C_{\mathcal{T}}$ are constraints on the corresponding variables.

3.1 State Transition

Consider a single model $\mathcal{M} = \mathcal{M}_i$ with corresponding transition operator $\phi = \phi_i$. For a state distribution $\mathbf{x}_k \sim \mathcal{N}(\boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k)$, the predicted observed state is $\mathbf{z}_k \sim \mathcal{N}(\mathbf{H}\boldsymbol{\mu}_k, \mathbf{H}\boldsymbol{\Sigma}_k\mathbf{H}^\top)$ and the measurement distributions is $\mathbf{y}_k \sim \mathcal{N}(\mathbf{H}\boldsymbol{\mu}_k, \mathbf{H}\boldsymbol{\Sigma}_k\mathbf{H}^\top + \boldsymbol{\Sigma}_y)$. Solving the optimisation problem in (6) requires the predictive distribution of the latent state \mathbf{x}_k . This means propagating the uncertainty in the inputs to the transition function f to its outputs. We assume we know *a priori* whether f is an analytical function or a black box, i.e. whether we *do* (analytical) or *do not* (black box) have derivative information of f with respect to its inputs. The derivative information is required for closed-form uncertainty propagation from inputs to outputs using Taylor approximations.

To obtain derivative information from a black-box transition function f , we place independent GP priors $f_{(d)} \sim \mathcal{GP}(m_{(d)}(\cdot), k_{x,(d)}(\cdot, \cdot)k_{u,(d)}(\cdot, \cdot)k_{\theta,(d)}(\cdot, \cdot))$ on output dimensions $d = 1, \dots, D_x$ of f . To simplify notation, let $\mu_f(\cdot) = \mathbb{E}_f[f(\cdot)]$ and $\Sigma_f(\cdot) = \mathbb{V}_f[f(\cdot)]$, such that

$$f(\cdot) \sim \mathcal{N}(\mu_f(\cdot), \Sigma_f(\cdot)) \equiv \begin{cases} \mathcal{N}(f(\cdot), \mathbf{0}), & f \text{ analytical} \\ \mathcal{N}(\mu(\cdot), \Sigma(\cdot)), & f \text{ black box} \end{cases} \quad (7)$$

where the posterior GP mean $\mu(\cdot)$ and covariance $\Sigma(\cdot)$ are given in (4).

Given an initial state estimate $\mathbf{x}_0 \sim \mathcal{N}(\boldsymbol{\mu}_0, \boldsymbol{\Sigma}_0)$, a sequence of control inputs $\mathbf{u}_k \sim \mathcal{N}(\hat{\mathbf{u}}_k, \boldsymbol{\Sigma}_u)$, $k = 0, \dots, T-1$, a model parameter posterior $\boldsymbol{\theta} \sim \mathcal{N}(\hat{\boldsymbol{\theta}}, \boldsymbol{\Sigma}_\theta)$, and the state transition described by (2), we wish to find the approximate state distribution $\mathbf{x}_k \sim \mathcal{N}(\boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k)$ at any time step $k \geq 1$, with mean and covariance given by the moments

$$\boldsymbol{\mu}_k = \mathbb{E}_{f, \mathbf{x}_0, \mathbf{u}_{0:t-1}, \boldsymbol{\theta}, \mathbf{w}_{0:t-1}} [\mathbf{x}_k], \quad (8a)$$

$$\Sigma_k = \mathbb{V}_{f, \mathbf{x}_0, \mathbf{u}_{0:t-1}, \boldsymbol{\theta}, \mathbf{w}_{0:t-1}}[\mathbf{x}_k]. \quad (8b)$$

Assume that the control covariance Σ_u , model parameter covariance Σ_θ and process noise covariance Σ_x are all constant and independent of \mathbf{x}_k and $\hat{\mathbf{u}}_k$.

Let $f(\tilde{\mathbf{x}}_k) \sim \mathcal{N}(\mu_f(\tilde{\mathbf{x}}_k), \Sigma_f(\tilde{\mathbf{x}}_k))$ denote the transition function evaluated at the concatenated state, control input and model parameters $\tilde{\mathbf{x}}_k = [\mathbf{x}_k^\top, \mathbf{u}_k^\top, \boldsymbol{\theta}^\top]^\top$, $\tilde{\mathbf{x}}_k \in \mathbb{R}^{D_x + D_u + D_\theta}$. Assuming the state, control input and model parameters are Gaussian distributed, the concatenated vector has Gaussian distribution $\tilde{\mathbf{x}}_k \sim \mathcal{N}(\tilde{\boldsymbol{\mu}}_k, \tilde{\Sigma}_k)$ with

$$\tilde{\boldsymbol{\mu}}_k = \begin{bmatrix} \boldsymbol{\mu}_k \\ \hat{\mathbf{u}}_k \\ \hat{\boldsymbol{\theta}} \end{bmatrix}, \quad \tilde{\Sigma}_k = \begin{bmatrix} \Sigma_k & \mathbf{0} & \text{cov}(\mathbf{x}_k, \boldsymbol{\theta}) \\ \mathbf{0} & \Sigma_u & \mathbf{0} \\ \text{cov}(\mathbf{x}_k, \boldsymbol{\theta})^\top & \mathbf{0} & \Sigma_\theta \end{bmatrix}, \quad (9)$$

where $\text{cov}(\mathbf{x}_0, \boldsymbol{\theta}) \equiv \mathbf{0}$, and $\text{cov}(\mathbf{x}_k, \mathbf{u}_k) \equiv \mathbf{0}$ (assuming $\mathbf{u}_k \neq \mathbf{u}_{k-1}$) since the latent state cannot depend on future control inputs.

To simplify notation, let $\nabla_{\boldsymbol{\xi}} g$, with $g \in \{f, \mu_f, \Sigma_f, \dots\}$, denote the partial derivative of $g(\boldsymbol{\xi})$ with respect to a variable $\boldsymbol{\xi}$, evaluated at the point $\mathbb{E}[\boldsymbol{\xi}]$.

The discrete-time state space model assumes process noise $\mathbf{w}_k \sim \mathcal{N}(\mathbf{0}, \Sigma_x)$. Using a first-order Taylor expansion of $\mu_f(\tilde{\mathbf{x}})$ around $\tilde{\mathbf{x}}_{k-1} = \tilde{\boldsymbol{\mu}}_{k-1}$, the mean and variance of the latent state at time step $k \geq 1$ in (8a) are approximately given by

$$\begin{aligned} \boldsymbol{\mu}_k &\approx \mu_f(\tilde{\boldsymbol{\mu}}_{k-1}), \\ \Sigma_k &\approx \nabla_{\tilde{\mathbf{x}}_{k-1}} \boldsymbol{\mu}_k \tilde{\Sigma}_{k-1} (\nabla_{\tilde{\mathbf{x}}_{k-1}} \boldsymbol{\mu}_k)^\top + \Sigma_x + \Sigma_f(\tilde{\boldsymbol{\mu}}_k), \\ \text{cov}(\mathbf{x}_k, \boldsymbol{\theta}) &\approx \nabla_{\boldsymbol{\theta}} \boldsymbol{\mu}_k \text{cov}(\mathbf{x}_{k-1}, \boldsymbol{\theta})^\top + \nabla_{\boldsymbol{\theta}} \boldsymbol{\mu}_k \Sigma_\theta. \end{aligned} \quad (10)$$

Note that $\nabla_{\tilde{\mathbf{x}}_{k-1}} \boldsymbol{\mu}_k \in \mathbb{R}^{D_x \times (D_x + D_u + D_\theta)}$. We calculate derivatives of $\boldsymbol{\mu}_k$ and Σ_k with respect to $\boldsymbol{\mu}_{k-1}$, Σ_{k-1} and $\hat{\mathbf{u}}_{k-1}$. The derivatives require second-order derivative information of f or the GP prediction. Appendix C and D have similar expressions to (10) for the latent state transition for continuous-time models and discrete-time models with Δ -transitions, respectively.

3.1.1 Combining Original Transition Function and GP Surrogate

The transition function f is replaced with a GP surrogate in (7) to enable analytical approximations of first- and second-order derivatives of f . Equation (7) presents a binary choice between an analytical approach where $\mu_f(\cdot) = f(\cdot)$, and a black-box approach where $\mu_f(\cdot) = \mu(\cdot)$. There is a third possible approach, where the black-box transition function f computes the predictive mean, and the GP surrogate only approximates the derivatives of f . Table 2 compactly shows the difference between the fully analytic approach (where f is analytic), the fully black-box approach (where only the GP surrogate is used) and the proposed third approach (where both f and the GP surrogate are used). The third approach limits the use of the GP surrogates to approximating the gradients of f .

This third approach is appropriate to use during model discrimination, i.e. when analysing agreement between model predictions and experimental observations. This reduces the risk that a model is discarded because of poor accuracy in the GP surrogate prediction. If f is expensive to evaluate, we may still choose to use the GP surrogate black-box approach to speed up design of experiments,

| Fully analytic approach | Fully black-box approach | Third approach |
|--|---|---|
| $\mu_f(\cdot) = f(\cdot),$ | $\mu_f(\cdot) = \mu(\cdot),$ | $\mu_f(\cdot) = f(\cdot),$ |
| $\Sigma_f(\cdot) = \mathbf{0},$ | $\Sigma_f(\cdot) = \Sigma(\cdot),$ | $\Sigma_f(\cdot) = \mathbf{0},$ |
| $\nabla\mu_f(\cdot) = \nabla f(\cdot)$ | $\nabla\mu_f(\cdot) = \nabla\mu(\cdot)$ | $\nabla\mu_f(\cdot) \approx \nabla\mu(\cdot)$ |

Table 2 When f is analytic we use the fully analytic approach in (7). When the transition function f is a black box we may choose to replace it in our computations with a GP surrogate for a fully black-box approach. A third approach is to use the original black-box transition function f for computing $\mu_f(\cdot)$ and the GP surrogate to approximate its gradients $\nabla\mu_f(\cdot)$.

i.e. when solving the optimisation problem in (6). Table 2 shows that for the fully analytic and black-box approaches, the gradient $\nabla\mu_f$ is exact, whereas for the third approach the gradient $\nabla\mu_f(\cdot) \approx \nabla\mu(\cdot)$ is an approximation. This may cause numerical issues when the GP mean $\mu(\cdot)$ —or rather the gradient $\nabla\mu(\cdot)$ —does not capture the behaviour in the transition function f with sufficient accuracy. If a numerical solver is provided with inaccurate gradients it may e.g. converge slowly to a solution, time-out before reaching any solution, or even throw an error. Therefore, from an optimisation point-of-view, it may be better to use the black-box approach than the third when solving the optimisation problem in (6). On the one hand this means solving an approximation of the optimisation problem, but on the other hand we may be more likely to find a solution.

3.2 Constraints

The optimisation problem in (6) is solved subject to constraints. Constraints are commonly related to safety concerns or physical limitations of the real systems, e.g. the maximum allowed electrical current in a machine or drug dose given to a patient. However, we may also need to apply constraints on the states $\mathbf{x}_{1:T}$ if the transition function f is replaced with a data-driven surrogate.

The initial state $\hat{\mathbf{x}}_0$, the control inputs $\hat{\mathbf{u}}_{1:T}$ and the measurement time points $\mathcal{T}_{\text{meas}}$ are independent, deterministic variables, set directly by the user. The states $\mathbf{x}_{1:T}$ and observed states $\mathbf{z}_{1:T}$ are dependent, stochastic variables. Constraints on independent and dependent variables are handled differently.

This section considers two types of constraints. Firstly, linear constraints

$$\mathbf{C}\boldsymbol{\xi} - \bar{\boldsymbol{\xi}} \geq \mathbf{0}, \quad (11)$$

where $\boldsymbol{\xi} \in \{\hat{\mathbf{x}}_0, \hat{\mathbf{u}}_k, \mathbf{x}_k, \mathbf{z}_k, t_k\}$, $\boldsymbol{\xi} \in \mathbb{R}^{D_\xi}$ is an independent or a dependent variable, $\mathbf{C} \in \mathbb{R}^{D_C \times D_\xi}$ and $\bar{\boldsymbol{\xi}} \in \mathbb{R}^{D_C}$, and the inequality is applied element-wise. The second type of constraints are constraints on the absolute difference, e.g. the rate of change, between independent variables. We will not consider constraints on the absolute difference between stochastic, dependent variables.

3.2.1 Independent Variable Constraints

The independent variables in the optimisation problem in (6) are the initial state settings $\hat{\mathbf{x}}_0$, the sequence of desired control inputs $\hat{\mathbf{u}}_{0:T-1}$ and the measurement

time points $\mathcal{T}_{\text{meas}}$ (for continuous-time models). Through Section 3.2.1, let $\boldsymbol{\xi} \in \{\hat{\boldsymbol{x}}_0, \hat{\boldsymbol{u}}_k, t_k\}$ denote an independent variable. Linear constraints on independent, deterministic variables are straight-forward to handle. Note that the constraint in (11) is written in the format of constraints in the (6) optimisation problem.

There may be limitations (for physical or safety reasons) to how quickly the control input $\hat{\boldsymbol{u}}_k$ can be varied. Let the absolute difference in dimension $d = 1, \dots, D_u$ of the control input between consecutive time steps be upper-bounded as $|\hat{u}_{k+1,(d)} - \hat{u}_{k,(d)}| \leq \Delta_{u,(d)}$. Using a standard reformulation, this constraint can equivalently be written as

$$\hat{u}_{k+1,(d)} - \hat{u}_{k,(d)} + \Delta_{u,(d)} \geq 0 \quad \wedge \quad \hat{u}_{k,(d)} - \hat{u}_{k+1,(d)} + \Delta_{u,(d)} \geq 0.$$

3.2.2 Dependent Variable Constraints

The dependent variables in the optimisation problem in (6) are the states $\boldsymbol{x}_{1:T}$ and observed states $\boldsymbol{z}_{1:T}$. Constraints on the dependent variables are typically more difficult to satisfy (Fu et al., 2015), because, as the name suggests, they are dependent on the initial state $\hat{\boldsymbol{x}}_0$ and the control sequence $\hat{\boldsymbol{u}}_{0:T-1}$.

Let model \mathcal{M} predict the state distribution $\boldsymbol{x}_k \sim \mathcal{N}(\boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k)$ and observed state distribution $\boldsymbol{z}_k \sim \mathcal{N}(\boldsymbol{\mu}_z, \boldsymbol{\Sigma}_z)$ at time step k , where $\boldsymbol{\mu}_z = \mathbf{H}\boldsymbol{\mu}_k$ and $\boldsymbol{\Sigma}_z = \mathbf{H}\boldsymbol{\Sigma}_k\mathbf{H}^\top$. Through Section 3.2.2, let $\boldsymbol{\xi}_k \sim \mathcal{N}(\boldsymbol{\mu}_\xi, \boldsymbol{\Sigma}_\xi) \in \{\boldsymbol{x}_k, \boldsymbol{z}_k\}$ denote a Gaussian-distributed dependent variable. We assume that the dependent variable constraint may be time-dependent, such that (11) becomes

$$\mathbf{C}_k \boldsymbol{\xi}_k - \bar{\boldsymbol{\xi}}_k \geq \mathbf{0}. \quad (12)$$

To simplify notation, let $\Xi_k \subset \mathbb{R}^{D_\xi}$ denote the space $\Xi_k = \{\boldsymbol{\xi} \mid \mathbf{C}_k \boldsymbol{\xi} - \bar{\boldsymbol{\xi}}_k \geq \mathbf{0}, \}$ at time step k , such that satisfying the linear constraint in (12) at time k is equivalent to satisfying $\boldsymbol{\xi}_k \in \Xi_k$. Multiple sources of uncertainty affect the states $\boldsymbol{\xi}_k$ and need to be accounted for. Constraints on dependent variables with unbounded probability distributions (e.g. Gaussian distributions) are referred to as *chance constraints*. The chance constraint equivalent of $\boldsymbol{\xi}_k \in \Xi_k$ is $P(\boldsymbol{\xi}_k \in \Xi_k) \geq 1 - \gamma$, for some $\gamma \in (0, 1)$ that determines the chance constraint's confidence level requirement. Chance constraints are typically analytically intractable (Prékopa, 1995, Ch. 11). There is a range of different tractable approximations for chance constraint, e.g. the scenario approach (Calafiore and Campi, 2006), the sample average approximation (Pagnoncelli et al., 2009), and the convex second-order cone approximation (Mesbah et al., 2014b). This paper only considers the cone approximation, and compares it to a dependent variable constraint that does not account for the uncertainty in the state $\boldsymbol{\xi}_k$.

Mean Constraint We let the term *mean constraint* denote the regular linear dependent variable constraint in (12) operating on the state mean, i.e. requiring $\boldsymbol{\mu}_\xi \in \Xi_k$. The mean constraint is used by existing literature on design of dynamic experiments (see Table 1). The mean constraint has the advantage of simplicity, at the expense of not accounting for the uncertainty in $\boldsymbol{\xi}_k$, and is approximately equivalent to solving the optimisation problem subject to $P(\boldsymbol{\xi}_k \in \Xi_k) \geq (\frac{1}{2})^{D_\xi}$. Hence, the mean constraint provides a poor guarantee that the constraint $\boldsymbol{\xi}_k \in \Xi_k$ will be satisfied for all time steps k .

Cone Constraint The convex second-order cone approximation (Mesbah et al., 2014b) decomposes the linear dependent variable constraint in (12) into multiple constraints $\mathbf{c}_{k,(j)}^\top \boldsymbol{\xi}_k - \bar{\xi}_{k,(j)} \geq 0$, for $\mathbf{C}_k = [\mathbf{c}_{k,(1)}, \dots, \mathbf{c}_{k,(D_C)}]^\top$ and $\bar{\boldsymbol{\xi}}_k = [\bar{\xi}_{k,(1)}, \dots, \bar{\xi}_{k,(D_C)}]^\top$. Each individual chance constraint $P(\mathbf{c}_{k,(j)}^\top \boldsymbol{\xi}_k - \bar{\xi}_{k,(j)} \geq 0) \geq 1 - \gamma$ can be satisfied by

$$\begin{bmatrix} \mathbf{c}_{k,(j)}^\top \\ \mathbf{c}_{k,(j)} \end{bmatrix} \boldsymbol{\mu}_\xi + \alpha \sqrt{\mathbf{c}_{k,(j)}^\top \boldsymbol{\Sigma}_\xi \mathbf{c}_{k,(j)}} \begin{bmatrix} 1 \\ -1 \end{bmatrix} - \begin{bmatrix} \bar{\xi}_{k,(j)} \\ \bar{\xi}_{k,(j)} \end{bmatrix} \geq \mathbf{0},$$

where $\alpha = \sqrt{2} \operatorname{erf}^{-1}(1 - \gamma)$, with $\operatorname{erf}^{-1}(\cdot)$ the inverse error function.

For a fixed covariance $\boldsymbol{\Sigma}_\xi$, the cone constraint is equivalent to the mean constraint for a smaller space $\hat{\Xi}_k \subset \Xi_k$. While mean constraint may poorly guarantee satisfying the chance constraint, the cone constraint may be overly conservative since it decomposes the full chance constraint into individual chance constraints.

3.2.3 State Constraints for Data-Driven Surrogate Models

When solving (6), the predicted states $\boldsymbol{\mu}_k$ may stray away from the state space region where there is state training data $\mathbf{X} = \{\mathbf{x}_1, \dots, \mathbf{x}_N\}$. This can cause numerical issues in the solver as the GP predictive variance grows large, and we may have reasons not to trust a corresponding allegedly optimal solution $\hat{\mathbf{u}}_{0:T-1}$. Hence $\boldsymbol{\mu}_k$ should be appropriately constrained as $\boldsymbol{\mu}_k \in \mathcal{X}$.

Assume we know the feasible control space \mathcal{U} defined by the input constraints. Control input training data can be sampled appropriately to fill the control space. We sample model parameter training data in a small region around the maximum *a posteriori* parameter estimate $\hat{\boldsymbol{\theta}}$ to compute estimated gradients for the Laplace approximation of $\boldsymbol{\Sigma}_\theta$. Assume the observed state $\mathbf{z}_k = \mathbf{H}\mathbf{x}_k$ is subject to a constraint $\mathbf{z}_k \in \mathcal{Z}, \forall k \in \mathcal{T}$. We would like to sample state training data from a domain $\mathcal{X}^* = \{\mathbf{x}_k \mid \exists \mathbf{u}_k \in \mathcal{U} : \mathbf{H}\mathbf{x}_k \in \mathcal{Z} \iff \mathbf{H}\mathbf{x}_{k+1} \in \mathcal{Z}\}$, i.e. states \mathbf{x}_k (i) that satisfy $\mathbf{z}_k \in \mathcal{Z}$, and (ii) for which some control input \mathbf{u}_k generates an observed state $\mathbf{z}_{k+1} \in \mathcal{Z}$. Finding \mathcal{X}^* is non-trivial even if the inverse transition function f^{-1} is known. \mathcal{X}^* can be approximated through exhaustive sampling, but this may be expensive. Proceeding, we let $\mathbb{Z}_x = \{1, \dots, D_x\}$ and assume $\mathcal{X}(\underline{\mathbf{x}}, \bar{\mathbf{x}}) = \{\mathbf{x} \mid \forall d \in \mathbb{Z}_x : \underline{x}_{(d)} \leq x_{(d)} \leq \bar{x}_{(d)}\} \approx \mathcal{X}^*$ is a known hypercube. GP training data is sampled from \mathcal{X} . We assume state training data is sampled from a grid (with samples on the borders of \mathcal{X}) and combined with control training data and model parameter samples. To maintain a fixed training data density, the training data set has to grow exponentially with state dimensionality D_x and control input dimensionality D_u .

For numerical stability when solving (6), we propose using state constraint $\boldsymbol{\mu}_k \in \mathcal{X}(\underline{\mathbf{x}} - \boldsymbol{\chi}_{d_{\text{out}}}, \bar{\mathbf{x}} + \boldsymbol{\chi}_{d_{\text{out}}})$ for the GP surrogate corresponding to output dimension d_{out} of f , $d_{\text{out}} \in \mathbb{Z}_x$, where $\boldsymbol{\chi}_{d_{\text{out}},(d)} \equiv |\bar{x}_{(d)} - \underline{x}_{(d)}| / \lambda_{(d_{\text{out}}),(d)}$ and $\lambda_{(d_{\text{out}}),(d)}$ denotes input dimension d 's lengthscale hyperparameter of output dimension d_{out} 's GP surrogate's state covariance function $k_{x,(d)}$. This ensures $\mathcal{X}(\underline{\mathbf{x}} - \boldsymbol{\chi}_{d_{\text{out}}}, \bar{\mathbf{x}} + \boldsymbol{\chi}_{d_{\text{out}}}) \subset \mathcal{X}(\underline{\mathbf{x}}, \bar{\mathbf{x}}) \subset \mathcal{X}(\underline{\mathbf{x}} - \boldsymbol{\chi}_{d_{\text{out}}}, \bar{\mathbf{x}} + \boldsymbol{\chi}_{d_{\text{out}}}), \forall d_{\text{out}} \in \mathbb{Z}_x$.

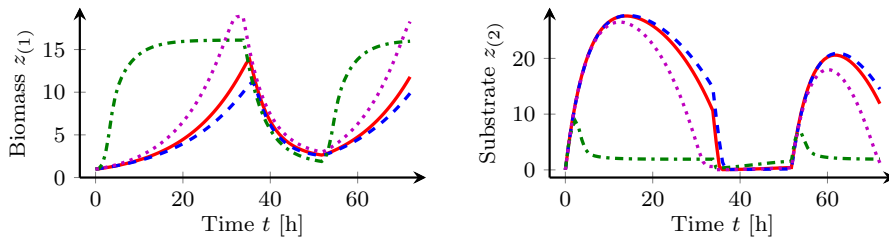


Fig. 2 Yeast fermentation case study model predictions for Espie and Macchietto (1989) optimal control inputs.

4 Results

We present computational results for an Espie and Macchietto (1989) case study that considers yeast fermentation, with constants, e.g. for true parameter values and noise covariances, from Chen and Asprey (2003). Yeast fermentation is common in pharmaceutical manufacturing (Martínez et al., 2012). There are $D_x = 2$ states (biomass and substrate concentration) and $D_u = 2$ control inputs. We observe both states, hence $D_z = D_y = 2$ and $\mathbf{H}_i = \mathbf{H} = \mathbf{I}$. The models have $D_{\theta,i} \in \{3, 4\}$ model parameters. Appendix E describes the case study in detail.

We present (i) a comparison to literature results for the Espie and Macchietto (1989) case study, (ii) a performance comparison of mean and cone constraints, (iii) simulations with correctly modelled, underestimated or overestimated control signal uncertainty, and (iv) a study of the GP surrogate performance as analytic emulators of black-box transition functions. We generated all results in our open-source Python package for design of experiments, *doepy*².

4.1 Comparison to Literature

Espie and Macchietto (1989) do not consider any uncertainty in their experimental design algorithm. Let \mathbf{U}_{EM}^* denote the optimal control signal found by Espie and Macchietto (1989). Fig. 2 shows the corresponding model predictions. We use the model parameter values θ_i reported in Espie and Macchietto (1989). If we perturb the elements of \mathbf{U}_{EM}^* by 1.5%, 5% and 10% and use the perturbed inputs as the initial point for the optimisation algorithm, we retrieve \mathbf{U}_{EM}^* again.

Let us add parametric uncertainty and a design criterion accounting for measurement noise. Chen and Asprey (2003) discriminate between the two most similar models: \mathcal{M}_1 and \mathcal{M}_2 . They find 20 optimal measurement time instances $\mathcal{T}_{\text{meas}}$ and an optimal control signal U_{CA}^* consisting of 5 piece-wise constant sections. Using the Chen and Asprey (2003) divergence measure, we find objective function value 2601 for their solution. For the same measurement time instances and control switch time instances, we find a solution $U_{\text{new},1}^*$ with objective function value 2805, i.e. a 7.8% larger divergence. We assume this discrepancy is due to the improvement in optimisation solvers after 20 years.

We add process noise with variance $\Sigma_x \equiv 0.01 \cdot \mathbf{I}$, i.e. the same order of magnitude as measurement noise. The optimal control signal U_{CA}^* of Chen and Asprey

² <https://github.com/scwolof/doepy>

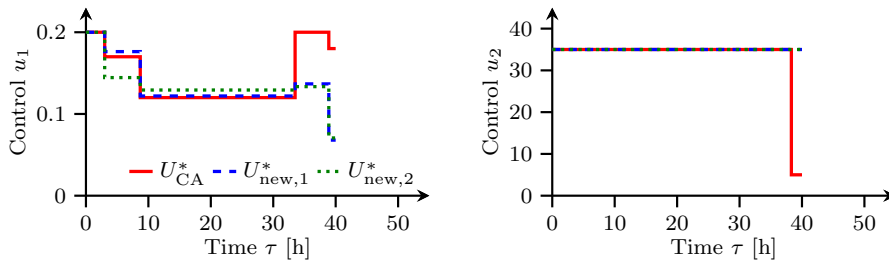


Fig. 3 U_{CA}^* from Chen and Asprey (2003), our solution $U_{new,1}^*$, and our solution $U_{new,2}^*$ for the case of added process noise.

(2003) yields a divergence of 244. Fig. 3 shows our solution $U_{new,2}^*$ that yields a divergence of 263, i.e. an 8% larger divergence. Fig. 3 shows that $U_{new,1}^* \neq U_{new,2}^*$, which means that adding process noise affects the optimal control signal.

4.2 Comparing Constraints

We compare the results of solving the design of experiments optimisation problem in (6) using (i) a mean constraint, or (ii) a cone chance constraint with $\alpha = 2$ standard deviations margin. We compare the constraints' performance in terms of the number of constraint violations, how severe the violations are, and the number of experiments required for successful model discrimination. This requires choosing a model discrimination criterion. The model discrimination literature describes several different discrimination criteria (Box and Hill, 1967; Buzzi-Ferraris et al., 1990; Michalik et al., 2010). Following Buzzi-Ferraris et al. (1990), we consider the χ^2 test. The weighted squared residuals $\delta_k^{(i)} = (\mathbf{y}_k - \mathbf{H}_i \boldsymbol{\mu}_k^{(i)})^\top (\mathbf{H}_i \boldsymbol{\Sigma}_k^{(i)} \mathbf{H}_i^\top + \boldsymbol{\Sigma}_y)^{-1} (\mathbf{y}_k - \mathbf{H}_i \boldsymbol{\mu}_k^{(i)})$ should be χ^2 -distributed. The χ^2 score is 1 minus the χ^2 cumulative distribution at $\sum_{k \in \mathcal{T}_{meas}} \delta_k^{(i)}$ with $|\mathcal{T}_{meas}| \times D_y - D_{\theta,i}$ degrees of freedom. Models are inadequate if their χ^2 score is below some threshold, e.g. $1E-3$.

Let \bar{z}_2 be a constraint upper bound on the substrate concentration, such that we wish to satisfy $z_{k,(2)} \leq \bar{z}_2$ for $k = \{1, \dots, T\}$. For the simulations, the upper bound takes a value $\bar{z}_2 \in \{7, 10, 15\}$. Measurements $\mathbf{y}_{1:T}$ are generated in each simulated experiment and used for model discrimination, with models deemed inadequate if their χ^2 score is below $1E-3$. Table 3 shows the performance in 25 simulations of the different state constraint types (mean constraint and cone constraint) in terms of the average number of experiments needed for successful model discrimination; All simulations correctly identify Model \mathcal{M}_1 as the true data-generating model. At least for the test instance by Chen and Asprey (2003), the additional conservatism of the cone constraints does not lead to (i) more required experiments or (ii) better model prediction performance.

But the following results indicate that the conservatism of the cone constraints do indeed increase the safety of the control signal by lowering the violation. We generate the experimental measurements $\mathbf{y}_{1:T}$ from one particular realisation of initial state, control inputs and measurement noise. If the constraints are not violated for this particular noise realisation, they may still be violated for other

| Bound | Constraint | Experiments required | | Avg. num. of models remaining after #n exp. | | |
|-------|------------|----------------------|------|---|------|----|
| | | Mean | Std | #1 | #2 | #3 |
| 7 | Mean | 2.02 | 0.14 | 2.00 | 1.02 | 1 |
| | Cone | 2.42 | 0.50 | 2.54 | 1.42 | 1 |
| 10 | Mean | 2.04 | 0.20 | 2.00 | 1.04 | 1 |
| | Cone | 2.16 | 0.37 | 2.16 | 1.16 | 1 |
| 15 | Mean | 2.12 | 0.33 | 2.02 | 1.12 | 1 |
| | Cone | 2.02 | 0.14 | 2.02 | 1.02 | 1 |
| 15 | Cone | 2.12 | 0.33 | 2.08 | 1.12 | 1 |

Table 3 Average number of experiments required (with standard deviation) to discard the incorrect models in 25 simulations of the yeast fermentation case study, with a mean or cone constraint (Section 3.2) with an upper bound on $z_{k,2}$. The right-most columns show the average number of models (out of four) that pass the χ^2 test after 1, 2 or 3 experiments. The last row shows the result for the GP surrogate approach in Section 4.4. For the test instance by Chen and Asprey (2003), the additional conservatism of the cone constraints does not lead to (i) more required experiments or (ii) better model prediction performance.

experiments generated using the same optimal control signal $\hat{\mathbf{u}}_{1:T}^*$. Therefore, we use Monte Carlo sampling to better assess the safety of a control signal.

Let \mathcal{M}_0 denote the data-generating model (\mathcal{M}_1 with “true” model parameter values). Assume we are studying an optimised control sequence $\hat{\mathbf{u}}_{0:T-1}^*$. We generate a large number $N_{\text{sim}} = 100$ of noisy control sequences $\mathbf{u}_{0:T-1,n}$, $n = 1, \dots, N_{\text{sim}}$ by drawing samples $\mathbf{u}_{k,n} \sim \mathcal{N}(\hat{\mathbf{u}}_k^*, \boldsymbol{\Sigma}_u)$. For each control sequence, we sample a random initial state $\mathbf{x}_{0,n} \sim \mathcal{N}([1, 0.01]^\top, \boldsymbol{\Sigma}_0)$ and process noise sequence $\mathbf{w}_{0:T-1,n}$. The control sequences, initial states and process noise sequences are used with the true model \mathcal{M}_0 to generate corresponding sequences of observed states $\mathbf{z}_{1:T,n}$. Let $\mathbf{Z} = \{\mathbf{z}_{1:T,1}, \dots, \mathbf{z}_{1:T,N_{\text{sim}}}\}$ denote the set of observed state sequences. Next, let $\mathbf{Z}_{\text{viol}} = \{\mathbf{z}_{1:T} \mid \mathbf{z}_{1:T} \in \mathbf{Z} \wedge \exists k : z_{k,(2)} > \bar{z}_2\}$ be the set of observation sequences for which the constraint $z_{k,(2)} \leq \bar{z}_2$ is violated for at least one time step k . The experiment *violation level* is the relative size $|\mathbf{Z}_{\text{viol}}|/N_{\text{sim}}$. A violation level of 0% means that a given control signal is apparently safe, while a violation level of 100% means that a control signal is almost guaranteed to result in constraint violations.

Fig. 4 shows the ratio of experiments, i.e. optimised control signals, at each violation level. Note: (i) The cone constraint results in *fewer* constraint violations than the mean constraint; (ii) The cone constraint results in *less severe* constraint violations than the mean constraint; (iii) The number of constraint violations increases as the upper bound shrinks. These observations are in line with expectations.

4.3 Effect of Control Input Uncertainty

Problem (6) assumes several covariance matrices are known. In many applications, the exact covariances are unknown and have to be approximated. We study the effect of under- or overestimating the covariance size by varying the control input covariance $\boldsymbol{\Sigma}_u$. More specifically, let the models \mathcal{M}_i assume the control input $\mathbf{u}_k \sim \mathcal{N}(\hat{\mathbf{u}}_k, \hat{\boldsymbol{\Sigma}}_u)$. We generate experimental data using control inputs $\mathbf{u}_k \sim \mathcal{N}(\hat{\mathbf{u}}_k, \boldsymbol{\Sigma}_u)$. In the simulations, the modelled covariance $\hat{\boldsymbol{\Sigma}}_u$ and the true covariance $\boldsymbol{\Sigma}_u$ are

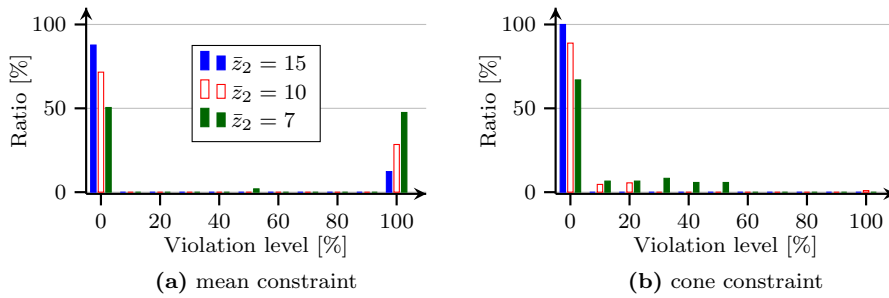


Fig. 4 The violation level is the ratio of Monte Carlo simulations using optimised control signals that result in constraint violations $z_{k,(2)} > \bar{z}_2$. The vertical axes show the ratio of optimised control signals for each violation level (horizontal axes) for a given constraint.

| Control covariance | | Experiments required | | Model discrimination | | | Cone constraint violations |
|--------------------|------------|----------------------|------|----------------------|-------|----------|----------------------------|
| $\hat{\Sigma}_u$ | Σ_u | Mean | Std | Succ. | Fail. | Inconcl. | |
| small | small | 2.04 | 0.20 | 100 % | 0 % | 0 % | 0 % |
| small | large | 2.05 | 0.22 | 84 % | 0 % | 16 % | 4 % |
| large | small | 2.58 | 0.97 | 96 % | 0 % | 4 % | 0 % |
| large | large | 2.21 | 0.66 | 96 % | 0 % | 4 % | 0 % |

Table 4 The first set of columns shows the modelled control covariance $\hat{\Sigma}_u$ and the true control covariance Σ_u used for generating experimental data, with $\Sigma_{\text{small}} = \text{diag}(1\text{E-}8, 1\text{E-}4)$ and $\Sigma_{\text{large}} = \text{diag}(1\text{E-}4, 1\text{E-}2)$. The second set of columns shows the average number of experiments required for successful model discrimination in 25 yeast fermentation case study simulations. The third set of columns shows the ratio of successful, failed or inconclusive model discrimination. The last column shows the ratio of simulations with cone constraint violations.

assigned values $\Sigma_{\text{small}} = \text{diag}(1\text{E-}8, 1\text{E-}4)$ and $\Sigma_{\text{large}} = \text{diag}(1\text{E-}4, 1\text{E-}2)$. There are four combinations of small and large modelled and true control covariances. The resulting scenarios can be described as (i) correctly modelled uncertainty $|\hat{\Sigma}_u| = |\Sigma_u|$, (ii) underestimated uncertainty $|\hat{\Sigma}_u| \leq |\Sigma_u|$, and (iii) overestimated uncertainty $|\hat{\Sigma}_u| \geq |\Sigma_u|$.

Table 4 shows the result of 25 simulations of the yeast fermentation case study with the different modelled and true control covariances. A cone constraint is enforced with an upper bound $\bar{z}_2 = 15$ (see Section 4.2). As expected, a correctly modelled small control covariance yields the best result in terms of average number of experiments required and the model discrimination success rate. A large modelled control covariance increases the average number of required experiments, and decreases the success rate. Model discrimination *fails* if an incorrect model is identified as the data-generating model, and is *inconclusive* if the experimental budget (maximum number of allowed experiments) is exhausted or the χ^2 -test discards all models as inadequate. These simulations never exhausted the experimental budget, i.e. all inconclusive model discrimination instances are due to all models being deemed inadequate. The rate of inconclusive model discrimination is significantly higher when the true control covariance is underestimated, and we have a cone constraint violation. Hence, we are punished less for conservative estimates of the control covariance than for overly optimistic estimates.

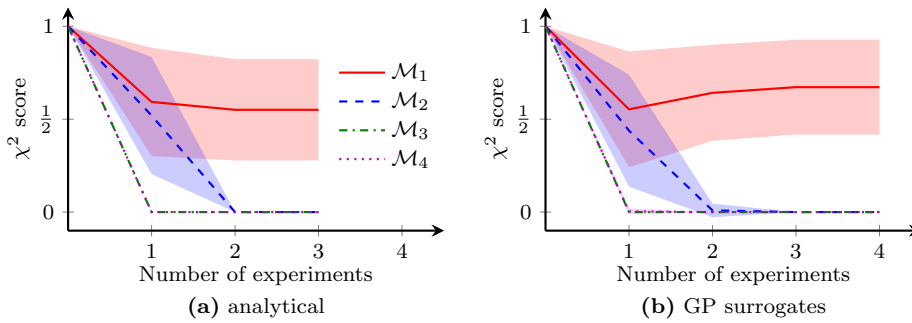


Fig. 5 Evolution of the average χ^2 scores (with one standard deviation) for the four rival models in the yeast fermentation case study, with experiments design using (a) the analytical approach, and (b) the GP surrogates approach.

4.4 Black-Box Transition Functions

Next we computationally compare the GP surrogate approach to the Section 4.2 analytic results. This approach uses GP surrogate predictions during design of experiments. Each simulation starts with no initial data and a relatively uninformed model parameter distribution. This model parameter distribution designs the first experiment. After executing an experiment, we update the model parameter distribution and then compute the models' χ^2 score during model discrimination. For the analytical transition functions in the Section 4.2 state constraint test and the Section 4.3 control uncertainty test, we use the Appendix E model parameter prior. However, for GP surrogates, which incorporate model prediction uncertainty, the model parameter prior is too large and the uncertainty cannot satisfy the constraint $\bar{z}_2 = 15$. We solve this by reducing the variance in the model parameter prior to $\Sigma_{\theta,i} = 1\text{E-}4 \cdot \mathbf{I}$ for the GP surrogates. On the one hand, this means the prior parameter estimate may be overly confident and there is a higher probability of constraint violations in the GP surrogate tests, but also the optimisation is more likely to converge on a feasible first experiment in each simulation.

The last row of Table 3 shows the GP surrogate approach performance in terms of average number of experiments required for successful model discrimination in 25 simulations of the yeast fermentation case study. Model \mathcal{M}_1 was correctly identified as the true data-generating model in all simulations. The GP surrogate approach has a marginally worse performance than the analytical method. Fig. 5 shows the average χ^2 score evolution (with one standard deviation) for the four rival models in the yeast fermentation case study. Fig. 5a shows the average χ^2 score for the analytical approach with a cone constraint, and Fig. 5b the average χ^2 score using the GP surrogates approach. This illustrates the GP surrogates' marginally worse performance compared to the analytical method. Black-box models yield larger marginal model prediction uncertainty that results in higher χ^2 scores.

5 Discussion

Our work assumes that the black-box model component is the state transition function f . Our problem formulation writes all models as first-order models, i.e.

| Approach | Mapping | Input dim. | Output dim. |
|------------------------|--|---------------------------------|----------------|
| Olofsson et al. (2018) | $\mathbf{x}_0, \mathbf{u}_{0:T-1}, \boldsymbol{\theta} \mapsto \mathbf{z}_{1:T}$ | $D_x + D_u \times T + D_\theta$ | $D_z \times T$ |
| This work | $\mathbf{x}_k, \mathbf{u}_k, \boldsymbol{\theta} \mapsto \mathbf{x}_{k+1}$ | $D_x + D_u + D_\theta$ | D_x |

Table 5 Comparison of the input and output dimensionality different GP surrogate models have for different mappings.

for discrete-time models the state \mathbf{x}_{k+1} depends only on the state \mathbf{x}_k , and for continuous-time models the first-order ordinary differential equation system have only the zeroth-order states $x(t)$ on the right-hand side. All n th order models can be expressed as first-order models by introducing additional states. We also assume the observed states \mathbf{z}_k are a linear combination of the states \mathbf{x}_k . Since the transition function f can be any non-linear function, using $\mathbf{z}_k = \mathbf{H}\mathbf{x}_k$ is only a minor restriction. State space models are also used in settings where the mapping from state to observed state is highly non-linear, e.g. in pixels-to-torque problems (Wahlström et al., 2015). Approximate moment matching can infer the predictive distribution for \mathbf{z}_k for this more general case $\mathbf{z}_k = g(\mathbf{x}_k)$.

Table 5 lists the input and output dimensionality of the GP mapping for both the approach in this paper and Olofsson et al. (2018). This approach would not work for design of dynamic experiments since (i) the input dimensionality would be too high for accurate GP inference and (ii) the output dimensionality would require many GP surrogate models. We reduce the input and output dimensionality by putting the GP prior on the state transition.

We consider an open-loop control approach, where a designed experiment runs to completion before data is analysed. Galvanin et al. (2009) and De-Luca et al. (2016) consider a closed-loop approach for parameter estimation. But in our setting, a closed-loop approach repeatedly solves the problem in (6). Although theoretically possible for designing dynamic experiments for model discrimination, the computational cost would likely be too high in practice.

6 Conclusions

We extended traditional analytic approaches for design of dynamic experiments using the Olofsson et al. (2018) methodology of replacing black-box models with GP surrogates. We also considered a wider range of problem uncertainty, and unified the literature contributions for discrete- and continuous-time models. Literature comparisons show that our models reduce to previously acquired results.

Both this paper and Olofsson et al. (2018) present the GP surrogate approach as hybridising analytical and data-driven approaches to optimal design of experiments. Condensing the range of possibilities, the GP surrogate approach is closer to the analytical than the fully data-driven approach. But we imagine a spectrum of different trade-offs between accuracy and computational complexity. The posteriors of more advanced GP models may not have closed form expressions (Salimbeni and Deisenroth, 2017), so we may require sampling to approximate predictive distributions. Hybrid design of experiments approaches using such advanced GP surrogates therefore lie closer to the original data-driven approaches. Model discrimination finds mechanistic models that adequately describe and predict a system’s behaviour. Mechanistic models are often needed in industry, e.g. to

satisfy regulatory requirements. In future, more models may be completely data-driven. An alternative, hybrid approach could combine mechanistic modelling and data-driven learning for data-driven models with physically meaningful embeddings (Sæmundsson et al., 2020). The model discrimination challenge becomes discerning the effect of the mechanistic versus the data-driven model part.

Conflict of interest. The authors declare that they have no conflict of interest.

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A Bania (2019) Models

The matrices \mathbf{A}_i , \mathbf{B}_i and \mathbf{C}_i in the Bania (2019) case study in (1) are defined as

$$\mathcal{M}_1: \quad \mathbf{A}_1 = -1, \quad \mathbf{B}_1 = 1, \quad \mathbf{C}_1 = 0.05, \quad (13)$$

$$\mathcal{M}_2: \quad \mathbf{A}_2 = \begin{bmatrix} 0 & 1 \\ -3 & -2.5 \end{bmatrix}, \quad \mathbf{B}_2 = \begin{bmatrix} 0 \\ 3 \end{bmatrix}, \quad \mathbf{C}_2 = \begin{bmatrix} 0 \\ 0.05 \end{bmatrix}, \quad (14)$$

$$\mathcal{M}_3: \quad \mathbf{A}_3 = \begin{bmatrix} 0 & 1 & 0 \\ -3 & -3.5 & 1 \\ 0 & 0 & -10 \end{bmatrix}, \quad \mathbf{B}_3 = \begin{bmatrix} 0 \\ 0 \\ 30 \end{bmatrix}, \quad \mathbf{C}_3 = \begin{bmatrix} 0 \\ 0 \\ 0.05 \end{bmatrix}. \quad (15)$$

The initial latent state is $\mathbf{x}_0^{(i)} = [0, \dots, 0]^\top$, and the process noise $w(t) \sim \mathcal{N}(0, 1)$ and the measurement noise $v_k \sim \mathcal{N}(0, 1)$ are Gaussian distributed.

B Analytical Design of Dynamic Experiments

Extensions of the analytical methods exist for design of dynamic experiments. Espie and Macchietto (1989) consider discrimination between multiple analytic continuous-time models and formulate an optimal control problem. Espie and Macchietto (1989) compare the results of using optimal constant control inputs *versus* an optimal dynamic control input. Chen and Asprey (2003) also consider continuous-time models and use a Laplace approximation for the model parameter covariance and linear propagation of the Gaussian model parameter uncertainty to approximate the marginal predictive distributions.

Skanda and Lebiedz (2010) assume Gaussian measurement noise to derive an expression for the Kullback-Leibler (KL) divergence between the predictive distributions of two rival models (with the same number of states). They include the measurement time points $\mathcal{T}_{\text{meas}}$ as variables in the optimisation problem, together with the initial state $x(0) = \mathbf{x}_0$ and control inputs $\mathbf{u}_{0:T-1}$. The control inputs Skanda and Lebiedz (2010) consider are additive perturbations to the state, and they assume the states cannot be measured and perturbed in the same time step. Skanda and Lebiedz (2013) extend the setup of Skanda and Lebiedz (2010) by considering model parameter uncertainty. They propose a robust optimisation formulation

$$\arg \max_{\substack{\mathcal{T}_{\text{meas}} \\ \mathbf{x}_0 \in \mathcal{X} \\ \mathbf{u}_{0:T-1} \in \mathcal{U}}} \min_{\substack{i, j \in \{1, \dots, M\} \\ i \neq j}} \min_{\substack{\boldsymbol{\theta}_i \in \boldsymbol{\Theta}_i \\ \boldsymbol{\theta}_j \in \boldsymbol{\Theta}_j}} \sum_{k \in \mathcal{T}_{\text{meas}}} \text{KL} [p(\mathbf{y}_k | \boldsymbol{\theta}_i) \| p(\mathbf{y}_k | \boldsymbol{\theta}_j)]$$

subject to constraints, with model parameter spaces $\boldsymbol{\Theta}_i$ and $\boldsymbol{\Theta}_j$ and $p(\mathbf{y}_k | \boldsymbol{\theta}_i)$ denoting the predictive distribution at time step k given model i with parameter values $\boldsymbol{\theta}_i$.

None of Espie and Macchietto (1989), Chen and Asprey (2003), or Skanda and Lebiedz (2010, 2013) consider process noise or uncertainty in the initial state $x(0) = \mathbf{x}_0$ or control signal $u(t)$. Nor do any of them, when solving the optimisation problem subject to dependent variable constraints on the observed states $\mathbf{z}_{1:T}$, account for the uncertainty in the observed states $\mathbf{z}_{1:T}$ predictions.

Cheong and Manchester (2014) consider non-parametric linear discrete-time systems with process noise (but no separate measurement noise) and uncertainty in the initial states \mathbf{x}_0 . For optimising the control signal they consider design criteria based on the pairwise difference in models' score in the χ^2 goodness-of-fit test. Though Cheong and Manchester (2014) consider dependent variable constraints in the observed states $\mathbf{z}_{1:T}$, they do not account for the uncertainty in the observed states $\mathbf{z}_{1:T}$ predictions.

Streif et al. (2014) and Mesbah et al. (2014a) look at cases of two rival non-linear models \mathcal{M}_1 and \mathcal{M}_2 with multiplicative measurement noise

$$\mathcal{M}_i: \quad \begin{cases} \left. \frac{dx(t)}{dt} \right|_t = f_i(x(t), u(t), \boldsymbol{\theta}_i), \\ z(t) = g_i(x(t), u(t), \boldsymbol{\theta}_i), \\ \mathbf{y}_k = \text{diag}(\mathbf{1} + \mathbf{w}_k) \mathbf{z}_k \end{cases}$$

where f_i and g_i , $i \in \{1, 2\}$, are polynomial functions. They consider uncertainty in the initial states and model parameters using polynomial chaos expansions, from which they compute higher moments of the predictive distributions—“a computationally formidable task” according to Streif *et al.* (2014). They discretise the control signal and solve the design problem by minimising the norm of the control signal such that the divergence between the predictive distributions is greater than or equal to some threshold value. The divergence can be computed using the predictive distributions’ higher moments (Streif *et al.*, 2014) or through Markov Chain Monte Carlo integration (Mesbah *et al.*, 2014a).

Keesman and Walter (2014) look at continuous-time models of the kind $\frac{d}{dt}y(t) = f(y(t)) + bu(t)$. They define the Hamiltonian and from this derive an optimal control law in closed form for two rival models. This requires gradient information of *at least* the first order. They do not account for parametric uncertainty, process noise or measurement noise.

Bania (2019) consider non-parametric linear discrete-time models with process and measurement noise. By looking at the mutual information between choice of model and observed output, they derive an optimisation formulation based on minimising the probability of selecting the wrong model. They mention how to extend their approach to non-linear models.

C Continuous-Time State Space Models

The continuous-time state space models is described by

$$\mathcal{M} : \begin{cases} \left. \frac{dx}{dt} \right|_t = f_i(x(t), u(t), \boldsymbol{\theta}) + w(t), \\ x(t_0) \equiv \mathbf{x}_0, \\ z(t) = \mathbf{H}x(t), \\ \mathbf{y}_k = z(t_k) + \mathbf{v}_k, \end{cases}$$

with state $x(t)$, control input $u(t)$, process noise distribution $w(t) \sim \mathcal{N}(\mathbf{0}, \boldsymbol{\Sigma}_x)$. The control input $u(t_k)$ may be a continuous function of the time t , but we will assume it is piece-wise constant. Model \mathcal{M} ’s state prediction at time step k is given by

$$\mathbf{x}_k = \mathbf{x}_{k-1} + \int_{t_{k-1}}^{t_k} f(x(t), \mathbf{u}_{k-1}, \boldsymbol{\theta}) dt + \mathbf{w}_k$$

where $w(t) \sim \mathcal{N}(\mathbf{0}, (t_k - t_{k-1})\boldsymbol{\Sigma}_x)$.

Let $\tilde{x}(t)$ denote the continuous concatenated state, control input and model parameters $\tilde{x}(t) = [x(t)^\top, u(t)^\top, \boldsymbol{\theta}^\top]^\top$, with Gaussian distribution $\tilde{x}(t) \sim \mathcal{N}(\tilde{\boldsymbol{\mu}}(t), \tilde{\boldsymbol{\Sigma}}(t))$, and let $\tilde{\boldsymbol{\mu}}_f$ denote the concatenated transition function

$$\tilde{\boldsymbol{\mu}}_f(t) = [\boldsymbol{\mu}_f(\tilde{\boldsymbol{\mu}}(t))^\top, \mathbf{0}, \mathbf{0}]^\top, \quad \tilde{\boldsymbol{\mu}}_f(t) \in \mathbb{R}^{D_x + D_u + D_\theta}.$$

We find the state prediction $\mathbf{x}_k \sim \mathcal{N}(\boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k)$ at time step k by extracting the corresponding elements from $\tilde{\boldsymbol{\mu}}(t_k)$ and $\tilde{\boldsymbol{\Sigma}}(t)$ (see (9)) which we compute by solving the following system of ordinary differential equations

$$\begin{cases} \left. \frac{d\tilde{\boldsymbol{\mu}}}{dt} \right|_t = \tilde{\boldsymbol{\mu}}_f(t), \\ \left. \frac{d\tilde{\boldsymbol{\Sigma}}}{dt} \right|_t = \nabla_{\tilde{\boldsymbol{\mu}}(t)} \tilde{\boldsymbol{\mu}}_f \tilde{\boldsymbol{\Sigma}}(t) + \tilde{\boldsymbol{\Sigma}}(t) (\nabla_{\tilde{\boldsymbol{\mu}}(t)} \tilde{\boldsymbol{\mu}}_f)^\top + \text{diag}(\boldsymbol{\Sigma}_f(\tilde{\boldsymbol{\mu}}(t)) + \boldsymbol{\Sigma}_x, \mathbf{0}, \mathbf{0}), \\ \tilde{\boldsymbol{\mu}}(t_{k-1}) \equiv \tilde{\boldsymbol{\mu}}_{k-1}, \\ \tilde{\boldsymbol{\Sigma}}(t_{k-1}) \equiv \tilde{\boldsymbol{\Sigma}}_{k-1}. \end{cases}$$

Derivatives of $\boldsymbol{\mu}_{k+1}$ and $\boldsymbol{\Sigma}_{k+1}$ with respect to $\boldsymbol{\mu}_k$, $\boldsymbol{\Sigma}_k$ and $\hat{\mathbf{u}}_k$ are calculated by integrating over the chain rule, and require second-order derivative information of f or the GP prediction.

For continuous-time models we may wish to optimise the measurement time instances $\mathcal{T}_{\text{meas}}$. In this case a minimum amount of time between measurements need to be enforced during optimisation. Let the absolute difference in time between two measurement time points t_k and $t_{k'}$ be lower-bounded by $\Delta_t \geq 0$

$$|t_k - t_{k'}| \geq \Delta_t, \quad \forall t_k, t_{k'} \in \mathcal{T}_{\text{meas}}. \quad (16)$$

This constraint is non-convex. To simplify the problem formulation, we introduce additional constraints to maintain a fixed order of the measurement time points, and reformulate (16) in convex form as

$$\forall t_k, t_{k'} \in \mathcal{T}_{\text{meas}} : \quad \begin{cases} t_k - t_{k'} - \Delta_t \geq 0, & k \geq k', \\ t_{k'} - t_k - \Delta_t \geq 0, & k < k'. \end{cases}$$

using the constraints format in (6).

D Discrete-Time Model with Δ -Transition

The discrete-time state space model with a Δ -transition is described by

$$\mathcal{M} : \quad \begin{cases} \mathbf{x}_{k+1} = \mathbf{x}_k + f(\mathbf{x}_k, \mathbf{u}_k, \boldsymbol{\theta}) + \mathbf{w}_k, \\ \mathbf{z}_k = \mathbf{H}\mathbf{x}_k, \\ \mathbf{y}_k = \mathbf{z}_k + \mathbf{v}_k, \end{cases}$$

with process noise $\mathbf{w}_k \sim \mathcal{N}(\mathbf{0}, \boldsymbol{\Sigma}_x)$. The discrete-time model with a Δ -transition follows from an Euler discretisation of continuous-time dynamics (Atkinson et al., 2009, Ch. 2). Using a first-order Taylor expansion of $\mu_f(\tilde{\mathbf{x}})$ around $\tilde{\mathbf{x}}_{k-1} = \tilde{\boldsymbol{\mu}}_{k-1}$, the mean and variance of the state at time step $k \geq 1$ in (8a) are approximately given by

$$\begin{aligned} \boldsymbol{\mu}_k &\approx \boldsymbol{\mu}_{k-1} + \mu_f(\tilde{\boldsymbol{\mu}}_{k-1}), \\ \boldsymbol{\Sigma}_k &\approx \nabla_{\tilde{\mathbf{x}}_{k-1}} \boldsymbol{\mu}_k \tilde{\boldsymbol{\Sigma}}_{k-1} \left(\nabla_{\tilde{\mathbf{x}}_{k-1}} \boldsymbol{\mu}_k \right)^\top + \boldsymbol{\Sigma}_x + \boldsymbol{\Sigma}_f(\tilde{\boldsymbol{\mu}}_k), \\ \text{cov}(\mathbf{x}_k, \boldsymbol{\theta}) &\approx \text{cov}(\mathbf{x}_{k-1}, \boldsymbol{\theta}) + \nabla_{\boldsymbol{\theta}} \boldsymbol{\mu}_k \text{cov}(\mathbf{x}_{k-1}, \boldsymbol{\theta})^\top + \nabla_{\boldsymbol{\theta}} \boldsymbol{\mu}_k \boldsymbol{\Sigma}_\theta. \end{aligned}$$

Note that $\nabla_{\mathbf{x}_{k-1}} \boldsymbol{\mu}_k = \mathbf{I} + \nabla_{\mathbf{x}_{k-1}} \mu_f$ with the Δ -transition model, and that $\nabla_{\tilde{\mathbf{x}}_{k-1}} \boldsymbol{\mu}_k \in \mathbb{R}^{D_x \times (D_x + D_u + D_\theta)}$. Derivatives of $\boldsymbol{\mu}_k$ and $\boldsymbol{\Sigma}_k$ with respect to $\boldsymbol{\mu}_{k-1}$, $\boldsymbol{\Sigma}_{k-1}$ and $\hat{\mathbf{u}}_{k-1}$ are calculated following the standard rules of matrix calculus, and requires second-order derivative information of f or the GP prediction.

It is common in GP regression to use zero-mean GP priors ($m_{(d)}(\cdot) \equiv 0$) to simplify calculations. The zero-mean prior is suitable for the Δ -transition state space model formulation (Ko et al., 2007; Deisenroth and Rasmussen, 2011).

E Yeast Fermentation Case Study

The yeast fermentation case study is taken from Espie and Macchietto (1989), with constants e.g. for true parameter values and noise covariances taken from Chen and Asprey (2003). There are $D_x = 2$ latent states (biomass and substrate concentration, respectively) and $D_u = 2$ control inputs (feed velocity and feed substrate concentration). We observe both states, hence $D_z = D_y = 2$ and $\mathbf{H}_i = \mathbf{H} = \mathbf{I}$, with $D_{\theta,i} \in \{3, 4\}$ model parameters. For simplicity, we omit the model index and time step index when writing out the models

$$\mathcal{M}_1 : \quad \begin{cases} \frac{dx_1}{dt} = (r - u_1 - \theta_4)x_1, \\ \frac{dx_2}{dt} = -\frac{rx_1}{\theta_3} + u_1(u_2 - x_2), \\ r = \frac{\theta_1 x_2}{\theta_2 + x_2}. \end{cases}$$

$$\mathcal{M}_2 : \begin{cases} \frac{dx_1}{dt} = (r - u_1 - \theta_4)x_1, \\ \frac{dx_2}{dt} = -\frac{rx_1}{\theta_3} + u_1(u_2 - x_2), \\ r = \frac{\theta_1 x_2}{\theta_2 x_1 + x_2}. \end{cases}$$

$$\mathcal{M}_3 : \begin{cases} \frac{dx_1}{dt} = (r - u_1 - \theta_3)x_1, \\ \frac{dx_2}{dt} = -\frac{rx_1}{\theta_2} + u_1(u_2 - x_2), \\ r = \theta_1 x_2. \end{cases}$$

$$\mathcal{M}_4 : \begin{cases} \frac{dx_1}{dt} = (r - u_1)x_1, \\ \frac{dx_2}{dt} = -\frac{rx_1}{\theta_3} + u_1(u_2 - x_2), \\ r = \frac{\theta_1 x_2}{\theta_2 + x_2}. \end{cases}$$

Data is generated from \mathcal{M}_1 with parameters $\theta = [0.25, 0.25, 0.88, 0.09]$ (Chen and Asprey, 2003). The measurement noise covariance is assumed known and given by

$$\Sigma_y = \begin{bmatrix} 0.06 & -0.01 \\ -0.01 & 0.04 \end{bmatrix},$$

and there is no process noise, hence $\Sigma_x \equiv \mathbf{0}$. We start without any experimental data and initial model parameter estimates $\theta_{i,d} = 0.5$ and covariance $\Sigma_{\theta,i} = 0.05 \cdot \mathbf{I}$ for all models \mathcal{M}_i and $d \in \{1, \dots, D_{\theta,i}\}$. The reason for this is the difficulty in finding initial experimental conditions that do not immediately render one or more models obviously inadequate. The initial states are given by $x_1 = 1$ and $x_2 = 0.01$, with initial latent state covariance $\Sigma_0 = \text{diag}(10^{-3}, 10^{-6})$. The controls have bounds $u_1 \in [0.05, 0.2]$ and $u_2 \in [5, 35]$. We let the control inputs have covariance given by $\Sigma_u = \text{diag}(10^{-6}, 10^{-3})$. We simulate 72 hours of fermentation, with measurements and changes in control signal every 1.5 hours.