Simulation-based Bayesian optimal design of aircraft trajectories for air traffic management

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SUMMARY

In this paper, we consider a specific Air Traffic Management problem, where multiple aircraft in a specific region are required to reach a different target zone in the minimum expected time, while maintaining safe separation. The problem is complicated by the presence of random wind disturbances. We propose a realistic policy to automatically generate optimal and safe maneuvers for each aircraft. The parameters of the optimal policy are computed using a sequential Monte Carlo approach. Copyright © 2010 John Wiley & Sons, Ltd.

Received 28 July 2009; Revised 3 July 2010; Accepted 7 July 2010

KEY WORDS: Bayesian optimal design; sequential Monte Carlo; air traffic management; conflict resolution

1. INTRODUCTION

In this paper, we address specific Air Traffic Management (ATM) problems, where multiple aircraft are interested in cruising through a monitored area to reach some specific point or area in the minimum expected time. The problem is complicated by the necessity of all aircraft maintaining safe separation. Although we will consider here only level flight, our proposed methodology is generic for a wide class of problems such as approaching terminal areas or changing flight levels. Our approach can be extended without much effort, for various aircraft dynamics, and for most established statistical wind models.

Planning aircraft trajectories to cross an area in the expected minimum time in the presence of constraints can be posed as a non-convex optimization problem. This is common in many problems in Air Traffic Management and the resulting costs or rewards in general often exhibit many local minima or maxima. This motivates the use of global optimization algorithms, which in some cases guarantee asymptotic convergence to a global optimum. In most cases such algorithms employ a randomized search strategy to ensure that the search process is not trapped in some local optimum. A popular example is Simulated Annealing (SA).

In addition to the issue of multiple local optima, solving such problems becomes more complicated when stochastic processes are used to represent model uncertainties, such as the effect of the wind velocity on each aircraft.

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The underlying distribution of the wind field is usually modeled by time-varying spatial stochastic processes. The resulting costs or rewards that need to be optimized are usually expectations over high-dimensional probability distributions, for which it seems impossible to devise analytical solutions or approximations except for a few very simple cases. Even computing such high-dimensional integrals can be a formidable task.

Markov Chain Monte Carlo (MCMC) methods [1], and Sequential Monte Carlo (SMC) methods [2–4] are currently the most successful methods for evaluating such expectations over relatively high-dimensional spaces under very weak assumptions, and have been widely applied in many areas such as finance, robotics and communications. An interesting point, often overlooked by the applied optimization community, is that Monte Carlo has also been applied for performing global optimization, mainly in inference problems such as Bayesian optimal design of experimental settings, standard Maximum Likelihood or Maximum a posteriori estimation [5–8]. Considerable theoretical support exists for both MCMC and SMC under very weak assumptions [1, 2].

Recently, in [9] an MCMC algorithm similar to SA developed in [8] has been applied for scheduling approaching aircraft to a terminal area and computing optimal way points for their maneuvers. In order to avoid the long execution times that are typically needed for stochastic optimization routines, it is important to consider carefully the computational complexity of the algorithms and take advantage of any structure that might be available in the problem to speed up computation. In the problems we will be considering this seems to be easier to manage when the SMC framework is used, yielding computationally more efficient algorithms, that are usually easier to implement. In addition, the implementation of SMC algorithms can be easily parallelized. In practice, this can lead to much shorter execution times, when these algorithms are implemented on conventional Graphics Processor Units (GPU) [10].

The outline of this paper is as follows: in Section 2 we formulate the problem as an open-loop control problem by constructing a dual-mode policy, which can be easily parameterized. In Section 3 we present a Bayesian interpretation of SA and in Section 4 we show how SMC can be used to compute the parameters of the optimal policy. In Section 5 we illustrate the performance of our approach using simulations. Finally, in Section 6 we provide some concluding remarks.

2. PROBLEM FORMULATION

In this section, we provide a probabilistic framework for a specific conflict resolution problem. We will consider level flight between $M$ aircraft cruising with a constant speed in a specified region. First, we will present the dynamics of each aircraft separately and then define the constraint region so that each aircraft maintains a safe separation from the rest, given some pre-specified tolerance probability. Then, we will consider the problem of finding optimal control inputs to meet these constraints, so that each aircraft arrives at a target exit region in the minimum expected time. Finally, we will use the following notation: for a $j$ dimensional vector we denote $z_{1:j} = (z_1, \ldots, z_j)$, $\|z\| = (z_1^2 + \cdots + z_j^2)^{1/2}$, $\delta_y(dz)$ the Dirac unit mass delta measure centered at $z = y$ and $\mathbb{I}_A(z)$ the indicator function of set $A$.

2.1. Dynamics

2.1.1. Dynamic model for each aircraft. We consider a standard two-dimensional Cartesian fixed-speed model for each aircraft $i$ [11], whose state $x_i^t$ ($\in \mathbb{R}^2$) is composed of the $x$ and $y$ Cartesian coordinates of the aircraft. Let the true aircraft speed of aircraft $i$ be $v_{\text{tas}}^i$ ($\in \mathbb{R}$), and let also the control input at each time be the heading, $h_i^t$ ($\in \mathbb{R}$, $\frac{\pi}{2}$), giving the following recursive dynamics for the state

\[
x_i^t = x_{i-1}^t + \begin{bmatrix} \cos h_i^t \\ \sin h_i^t \end{bmatrix} v_{\text{tas}}^i \delta + w_{t-1}(x_{i-1}^t) \delta
\] (1)
where \( \delta \) is the common time discretization interval, and \( w_t(x) \) is a random disturbance representing the wind velocity at point \( x \) of the state space at time \( t \). We assume that the sequence of headings \( h_{1:t}^i \) can be parameterized by some fixed (possibly low) dimensional parameter \( \theta^i \) and \( w_{t-1} \) is a spatio-temporal varying random variable defined on some open subset of \( \mathbb{R}^2 \). Also, for each aircraft \( i \), we will assume that there exists some exit point \( e^i \) (\( \in \mathbb{R}^2 \)). The aim of the control inputs is to drive each aircraft \( i \) to some neighborhood of \( e^i \), denoted by \( B^i (\subset \mathbb{R}^2) \), which we will refer to frequently in this paper as the ‘target set’ or ‘exit region’.

### 2.1.2. Parameterization of control inputs

We will show how \( h_{1:t}^i \) can be represented as a deterministic function of a parameter \( \theta^i \). Our approach uses ‘blocking’ of control inputs to reduce the dimension of the control space. The motivation for doing this arises from both the practical issues in the specific application and the computational aspects of the underlying methodology. As far as the application is concerned the resulting maneuvers need to be simple enough to be implemented in practice either by a pilot or an autopilot. In terms of the methodology it is also important to reduce the control space to a reasonable size parameter space, so that a fast and efficient search for the optimal control law can be made possible. We will refer to the control law also as a ‘policy’ and consider the parameter of the policy to be a deterministic variable to be computed.

More specifically for our parameterization, we will divide the sequence of \( h_{1:t}^i \) into a sequence of \( H \) blocks of constant inputs as follows:

\[
h_{1:t}^i = \left[ [h_1^i, \ldots, h_1^i], [h_2^i, \ldots, h_2^i], \ldots, [h_H^i, \ldots, h_H^i] \right],
\]

Block 1 Block 2 Block H

In this case, for every \( k = 1, \ldots, H \), each \( h_k^i \) is composed of the stacked vector of every \( h_k^i \), which in turn can be computed by using the change in heading, namely \( \theta_k^i \), compared with the heading in the first block \( h_1^i \). This gives the following control update for \( k > 1 \):

\[
h_k^i = h_1^i + \theta_k^i,
\]

where we will for the time being set \( \theta_1^i = 0 \). Then the parameter is simply given by \( \theta^i = [\theta_1^i, \theta_2^i, \ldots, \theta_H^i] \). Note that we have not so far made any assumptions on each block size, which may vary with \( k \) and could be used as a design or tuning variable.

In this paper we shall use the same length \( p \), for each block, except the last one, which will usually be a longer block. In the last block we propose the use of either **tracking a straight-line course** or a **fixed heading**:

**Tracking:** We propose that each element in the final block is time varying in this case. In tracking the heading counteracts the cross-track action of the wind [12]. The inputs are computed at time \( t \) so that \( x_{t+1}^i \) tracks the straight line connecting \( x_{t-(H-1)p}^i \) to the exit point \( e^i \) for aircraft \( i \). Therefore, each element in \( h_{H}^i = [h_{(H-1)p+1}^i, \ldots, h_H^i]^T \) can be computed recursively for \( m = t - (H - 1)p + 1, \ldots, t \) as follows:

\[
h_{m+1}^i = \varphi_m^i + \sin^{-1} \left( \frac{\| w_m(x_m^i) \| \sin \left( \tan^{-1} \left( \frac{w_m(x_m^i)(2)}{w_m(x_m^i)(1)} - \varphi_m^i \right) \right)}{\delta v_{\text{as}}} \right)
\]

where $\| \cdot \|$ denotes the 2-norm and for any time $t$ and aircraft $i$, the current angle, $\varphi_i$, toward the exit-point $e^i$, is given by

$$\varphi_i = \tan^{-1}\left( \frac{e^i(2) - x^i(2)}{e^i(1) - x^i(1)} \right)$$

(4)

**Fixed heading:** In this case we propose that $h^i_H$ is the stacked vector composed of $t - (H - 1) p$ repetitions of $h^i_{H-1}$, which is given by the heading angle after the last input of $h^i_{H-1}$ is applied, i.e.

$$h^i_H = \varphi^i_{t-(H-1)p}$$

(5)

Note that in both cases, in contrast to the previous blocks, the size of the final block is not fixed but varies depending on $t$. In our application $t$ will keep growing until $x^i_t$ has reached its target set $B^i$. As a result the size of the last block can be thought of as random, but in each case the input heading is given by a deterministic update scheme, as shown in (3)-(5). The benefit of this strategy in each case is that $\theta^i$ is fixed and possibly low-dimensional, regardless of the time required for the state trajectories to reach the target sets. The final block is somewhat analogous to the ‘terminal control law’ that is commonly used in ‘dual-mode’ model predictive control schemes, as introduced in [13]—we think of the first $H-1$ blocks as ‘mode 1’, and of the final block as ‘mode 2’.

### 2.1.3. Multiple aircraft state augmentation and constraints.

So far we have considered the dynamics of each aircraft separately. Now, we will assume that $M$ aircraft are always present in the area of surveillance, so $i = 1, \ldots, M$. Also we denote $x_i$ to be the stacked vector of every $x^i$, i.e. $x_i = [x^i_1, \ldots, x^i_M]^{\top}$ and similarly for the control parameters we have $\theta = [\theta^1, \ldots, \theta^M]^{\top}$.

We will also assume that for a fixed known value of $\theta$, $X_t$ is a Markov chain generated by the law $P_{x_0,\theta,t}(dx_{0:t}) = \prod_{k=1}^{t} f(x_k| x_{k-1})dx_{0:k}$ with initial distribution $\delta_{x_0}(dx)$. The dynamics can be also written as a joint Markov transition density $f_\theta$ with respect to the Lebesgue measure $dx_t$:

$$P_{x_0,\theta}(X_t \in dx_t | X_{1:t-1} = x_{1:t-1}) = P_{x_0,\theta}(X_t \in dx_t | X_{t-1} = x_{t-1}) = f_\theta(x_t | x_{t-1})dx_t.$$  

(6)

Each individual transition density $f_\theta(x^i_t | x_{t-1})$ with respect to $dx^i_t$ defined as

$$P_{x_0,\theta}(X^i_t \in dx^i_t | X^i_{1:t-1} = x^i_{1:t-1}) = P_{x_0,\theta}(X^i_t \in dx^i_t | X^i_{t-1} = x^i_{t-1}) = f_\theta^i(x^i_t | x_{t-1})dx^i_t$$

can be obtained through marginalization

$$f_\theta^i(x^i_t | x_{t-1}) = \int f_\theta(x_t | x_{t-1})dx_{-i}^t,$$  

(7)

where $x^{-i} = [x_1, \ldots, x^i-1, x^{i+1}, \ldots, x^M]$.

For the purpose of ATM, we are primarily interested in maintaining at least some minimum safe separation distance between each aircraft. Therefore, at each time $t$ for the state vector $x_t$ we define the feasible state space as

$$\mathcal{X}_t = \{x^i_t \in \mathbb{R}^2 : \forall j \neq i \|x^j_t - x^i_t\| > d, i, j = 1, \ldots, M\},$$

(8)

and the path space of the feasible trajectories as $\mathcal{F}_0 = \mathcal{X}_0 \times \cdots \times \mathcal{X}_t$. Moreover, denote $\Theta$ as the set of admissible actions $\theta$ such that for a given $t$ the probability of safe separation ($\text{PSS}_t$) satisfies

$$\text{PSS}_t := P_{x_0,\theta,t}(X_{0:t} \in \mathcal{F}_{0:t}) \geq 1 - \epsilon,$$  

(9)
where $\varepsilon \in [0, 1)$ is some tolerance level. Note that in problems with stochastic dynamics it might not be possible to use deterministic constraints and we therefore opt for the use of some pre-specified tolerance level $\varepsilon$ which can be interpreted as an allowed or safe probability of conflict. Note that, although it was natural to pose the joint dynamics in a distributed fashion as done for the case of each aircraft earlier, this is not straightforward for the case of the constraints, which involve pair-wise comparisons of the aircraft’s states.

2.1.4. Feynman–Kac model for the augmented process. In order to avoid dealing with the random sets $\mathcal{X}_t$ defined in (8) we will use the discrete-time Feynman–Kac model formalism described in detail in [2]. For convenience, let $X_t = X_{1:t}$ denote the sequence of $t$ consecutive augmented states which is generated by the law of the Markov chain $P_{x_0, t, t}(dx_{0:t})$. We will use the potential function $G_k(x_k) = \prod_{k=1}^{t-1} G_k(x_k)$ where

$$G_k(x_k) = \begin{cases} 1 & \inf_{i \neq j} \| x_k^i - x_k^j \| > d, \quad i, j = 1, \ldots, M \\ 0 & \text{otherwise} \end{cases}$$

to define the probability measure of the path of feasible state sequences up to time $t$ as

$$Q_{x_0, t, t}(X_t \in \mathcal{A}) = \frac{1}{Z_t} G_t(X_t) Q_{x_0, t, t}(X_t \in \mathcal{A}),$$

where $\mathcal{A} \in \mathbb{R}^{2M}$ is an arbitrary set and $Z_t$ is the normalizing constant given by

$$Z_t = \mathbb{E}_{P_{x_0, t, t}} \left[ \prod_{k=1}^{t-1} G_k(x_k) \right],$$

with $\mathbb{E}_{P_{x_0, t, t}}$ denoting the expectation under the probability measure $P_{x_0, t, t}$.

An alternative definition for the law of the feasible state sequence path up to time $t$ can be given using an arbitrarily chosen bounded test function $\Phi : \mathbb{R}^{2t} \rightarrow \mathbb{R}$ as follows:

$$Q_{x_0, t, t}(\Phi) = \frac{1}{Z_t} \mathbb{E}_{P_{x_0, t, t}} \left[ \Phi(X_{0:t}) \prod_{k=1}^{t-1} G_k(x_k) \right].$$

Similarly, using a bounded test function $\phi : \mathbb{R}^2 \rightarrow \mathbb{R}$ we can derive a sequence of distributions $\eta_t$ for each $X_t$ as

$$\eta_t(\phi) := \frac{\eta_t(\phi)}{\gamma_t(1)} = \frac{1}{Z_t} \mathbb{E}_{P_{x_0, t, t}} \left[ \phi(X_t) \prod_{k=1}^{t-1} G_k(x_k) \right].$$

(10)

Given $\eta_k(G_k) = Z_{k+1}/Z_k$ it is straightforward to show that for the unnormalized measure $\gamma_t$ that

$$\gamma_t(\phi) = \eta_t(\phi) \prod_{k=1}^{t-1} \eta_k(G_k).$$

(11)

For the problem of interest one can interpret the normalizing constant $\gamma_t(G_t) = \gamma_{t+1}(1)$ as being $PSS_t$. Consequently from (10)–(11) we obtain the recursion

$$PSS_t = \gamma_{t+1}(1) = Z_{t+1} = \prod_{k=1}^{t} \eta_k(G_k),$$

(12)

which will be used later to compute $PSS_t$ numerically.
2.1.5. Monte Carlo approximations and relaxing the constraints. Although the Feynman–Kac model described earlier provides a convenient interpretation of PSS\(_t\) as the normalizing constant \(Z_{t+1}\), we note that in most cases this normalizing constant cannot be computed analytically. In this section we will briefly sketch how to sample from \(Q_{x_0,0,t}\) and numerically estimate PSS\(_t\) by a Monte Carlo approximation based on Importance Sampling (IS) [4]. The aim is to relax the constraints, so that they remain valid for the approximation of PSS\(_t\) instead.

First we assume that it is possible to sample from the true dynamics using samples from the transition density \(f_0\) and we can therefore obtain \(\kappa\) independent identically distributed (iid) samples \(\{x_t^{(i)}\}_{i=1}^{\kappa}\) from \(\mathbb{P}_{x_0,0,t}\). We can then check whether any conflicts are present at any time and select samples as samples of \(Q_{x_0,0,t}\) only when \(\mathcal{G}_t(x_t^{(i)}) = 1\). The Monte Carlo estimate of PSS\(_t\), namely \(\hat{\text{PSS}}_t\), is given by:

\[
\hat{\text{PSS}}_t = \frac{1}{\kappa} \sum_{i=1}^{\kappa} \mathcal{G}_t(x_t^{(i)})
\]

One could then soften the constraints to take the following form, \(\hat{\text{PSS}}_t \geq 1 - \varepsilon\).

In the subsequent methodology we will assume that it is indeed possible to simulate from the true transition density \(f_0\). Even if this is not the case our methodology can still be used provided that one can still obtain unbiased samples \(x_1:t\) from the dynamics. The minimal requirement is that point-wise computation of \(f_0\) is available. Therefore, the methodology and algorithms presented in the remainder of this paper are generic and can be used not only for level flight, but also for example in problems with changing altitudes.

2.2. Cost function

We define for each aircraft the time of arrival \(T^i\) at the target set \(B^i\) as

\[
T^i = \inf_{t \geq 1} \{ t : X^i_t \in B^i \}.
\]

Note that in the context of stochastic processes the time of arrival corresponds to the random hitting time of set \(B^i\). We will assume that once some aircraft first enters into \(B^i\) with a realization of the first hitting time \(T^i\) being \(\tau^i\), then the state corresponding sequence \(x_{\tau^i+m}^i\), \(m > 1\), remains in an arbitrary absorbing feasible set, where the subsequent trajectory does not violate the constraints.

We will write the expected time of arrival or hitting time at \(B^i\) as:

\[
\mathbb{E}_{x_0,0}[T^i] = \sum_{\tau=1}^{\infty} \tau^i \pi_{T^i,0}(\tau^i) \text{ with } \pi_{T^i,0}(\tau^i) = \frac{Q_{x_0,0,\tau^i+1}\{T^i = \tau^i\}}{\sum_{m=1}^{\infty} Q_{x_0,0,m+1}\{T^i = m\}}. \tag{13}
\]

We define the probability mass function (pmf) of the hitting times (or times of arrival), \(\pi_{T^i,0}(\tau^i) \propto Q_{x_0,0,\tau^i+1}\{T^i = \tau^i\}\), using \(Q_{x_0,0,t}\) instead of the natural law of the process \(\mathbb{P}_{x_0,0,t}\) to avoid \(G_k(X_k) = 0\) occurring, which would lead to an infinite \(T^i\). Thus we ensure that for all times \(t \leq \tau^i\), the trajectories \(X^i_t\) remain in the feasible space \(\mathcal{X}_{0:t}\). Furthermore, to avoid other pathological cases, we will also assume that for every \(i\), \(Q_{x_0,0}\{T^i < \infty\} = 1\) holds for some value of \(\theta \in \Theta\) and in addition \(\mathbb{E}_{x_0,0}(T^i) < \infty\) for every \(\theta, i\).

In order to compute \(\pi_{T^i,0}\) we will consider all possible feasible state sequences \(X^i_{\tau^i}\), for which the state of aircraft \(i\) just reaches the target set \(B^i\) at time \(\tau^i\) provided the previous part of the trajectory has not done so. For this reason we define the function \(g^i_{\tau^i} : \mathbb{R}^{2M \tau^i} \rightarrow \{0, 1\}\)

\[
g^i_{\tau^i}(x_{\tau^i}) = \prod_{k=1}^{\tau^i-1} \mathbb{1}_{x^i_k \not\in B^i} \prod_{k=1}^{\tau^i-1} \mathbb{1}_{x^i_k \not\in B^i}. \tag{14}
\]
and use it to compute \( \pi_{T^i, \theta}(\tau^i) \) as
\[
\pi_{T^i, \theta}(\tau^i) \propto Q_{x_0, \theta, \tau^i + 1} \{ T^i = \tau^i \} = \frac{\mathbb{E}_{p_{x_0, \theta, \tau^i}} \left[ g_{T^i}^i(x_{\tau^i}) \prod_{k=1}^{\tau^i} G_k(x_k) \right]}{\mathbb{E}_{p_{x_0, \theta, \tau^i}} \left[ \prod_{k=1}^{\tau^i} G_k(x_k) \right]} = \eta_{\tau^i + 1}(g^i) \tag{15}
\]

It is implied that by sampling any trajectory from \( Q_{x_0, \theta, \tau^i + 1} \), one can determine if \( \tau^i \) is a sample from \( \pi_{T^i, \theta}(\tau^i) \) by checking whether \( g_{\tau^i}^i(x_{\tau^i}) = 1 \). In this case we will require that for every time \( t > \tau^i \), \( g^i_t \) is also holds, because \( B^i \) was formulated as an absorbing set.

In this paper we shall consider minimizing the expected maximum time of arrival with respect to all aircraft. More precisely we define the stopping time to be minimized as \( T = \max_i T^i \) and cast the problem so as to compute
\[
\theta^* = \arg \min_{\theta \in \Theta} \mathbb{E}_{x_0, \theta}[T] \quad \text{s.t. } \mathbb{P}_{x_0} \geq 1 - \varepsilon
\]
where the relaxed constraints impose the need of all aircraft staying within the feasible region with some confidence probability \( 1 - \varepsilon \). This is a difficult cost to minimize by conventional methods, such as analytical computations or the use of gradients. This motivates the use of simulation-based methods to estimate \( \theta^* \).

In addition, we can use that \( \mathbb{E}_{x_0, \theta}(T) < \infty \) and therefore can be bounded from above by a constant \( \tau_{x_0} = \sup_{\theta} \mathbb{E}_{x_0, \theta}(T) \), to transform the minimization problem into a maximization problem and estimate \( \theta^* \) as the maximizers of
\[
J_{x_0}(\theta) = \int u(\tau) \pi_{T, \theta}(\tau) \, d\tau, \tag{16}
\]
where
\[
u(\tau) = \tau - \tau \quad \text{and} \quad \pi_{T, \theta}(\tau) = \sum_{t=1}^{\infty} \sum_{m=1}^{\infty} \mathbb{Q}_{x_0, \theta, \tau, m + 1} \{ T = \tau \} \mathbb{Q}_{x_0, \theta, m + 1} \{ T = m \} \delta_{\tau}(d\tau) \tag{17}
\]
and \( \tau > \tau_{x_0} \) is an arbitrary constant so that \( u > 0 \). Having to solve a maximization problem makes it easier to use the annealing methods found in [5–8].

Note that at first glance choosing such a \( \tau \) might seem restrictive, but this is mitigated by the fact that in practice it might be easy to find some \( \tau \) large enough so that the problem remains unaffected. To simplify exposition in the subsequent sections we formulated \( \pi_{T, \theta}(\tau) \) as a probability density function composed from a train of Dirac delta functions. As regards to computing \( \pi_{T, \theta}(\tau) \) one can use \( \mathbb{Q}_{x_0, \theta, \tau + 1} \{ T = \tau \} \propto \eta_{\tau + 1}(\prod_{i=1}^{M} g_{i}^i) \) similar to Equation (15).

### 2.3. Wind modeling

In this section we shall discuss how the wind velocity \( w_t(x) \) can be modeled at each point \( x \in \mathbb{R}^2 \). Our approach is close to that of [14–16] and similar ideas have appeared in [17]. More formally, we assume that \( W_t \) is a measurable process taking values in \( \mathbb{R}^2 \). We will assume that the mean of the process \( \mu_t \) is a spatio-temporal varying two-dimensional surface, that depends on the past history of the process through a Markov relation
\[
\mathbb{E}[W_t(x)|W_{t-1} = w_{t-1}] = \mathbb{E}[W_{t-1}(x)|W_{t-1} = w_{t-1}] = \mu_t(x; s_{\mu}(w_{t-1})),
\]
where \( s_{\mu}(w_{t-1}) \) is a vector of fixed dimensional sufficient statistics, for example some linear combination of \( s_{\mu}(w_{t-2}) \), the previous mean and covariance. Similarly for the variance we will assume
\[
\mathbb{E}[(W_t(x) - \mu_t(x; s_{\mu}(w_{t-1})))^2|W_{t-1} = w_{t-1}] = \Sigma_t(x; s_{\Sigma}(w_{t-1})),
\]
with $s_\Sigma(w_{t-1})$ being a sufficient statistics for the variance. As we are assuming a stationary model, note that the sufficient statistics can be computed recursively.

This level of abstraction formulates the problem in its most general setting, but for the purposes of this paper, we will use the linear Gaussian model

$$w_t(x) = \mu_t(x; s_\mu(w_{t-1})) + v_t,$$

where $v_t \sim N(0, \Sigma_t(x; s_\Sigma(w_{t-1})))$. In [14–16] the authors assumed that $s_\mu$ and $s_\Sigma$ depend on the complete path of $W_{1:t-1}$. This choice was influenced by the exhaustive data analysis on weather forecasts carried out in [18]. There the data sets analyzed spanned a very large period of time, order of magnitudes larger than the few tens of minutes we are interested in. In this paper we choose to use Markov modeling rather than using a time-varying autoregressive process as in [19]. This was further validated in the data sets used for simulations in Section 5. Again, we stress that the methodology and the simulation-based algorithm presented later in this paper does not rely on this specific model choice.

In practice, the surface $\mu_t(x)$ is unknown and what is available is gridded data, which corresponds to Rapid Update Cycle (RUC) weather forecasts [18] at specific points of the state space. The forecasts typically appear in intervals of three hours for the next three days. They are updated daily and can be downloaded from the National Oceanic and Atmospheric Administration (NOAA) server. We shall use this gridded data to provide smooth approximations of $\mu_t$ and $\Sigma_t$ both in space and time.

As far as $\mu_t$ is concerned, obtaining approximations for the surfaces drops to a regression problem. For both spatial and temporal regression we opted for the use of plain linear interpolation as in [16]. These choices were largely influenced by the inherent simplicity and computational efficiency. In addition early, but not exhaustive, numerical validation indicated that using more advanced regression techniques, such as Bayesian analysis [20], or splines [17] for the spatial regression, could yield better approximations but at a significant computational cost. In any case, the benefit of these methods is marginal in terms of illustrating the performance of our approach. Therefore these and other advanced methods should not be ruled out, when it comes to implementing a more realistic or complex example, but their further investigation is beyond the scope of this paper.

As regards obtaining approximations for $\Sigma_t$, we used a simplification of the work found in [14, 16], suitable for implementation of the spatial component of a Markovian wind field. Based on the additional assumption that the variance retains a diagonal isotropic structure for the two-dimensional model, we use the following covariance:

$$\mathbb{E}[(W_t - \mu_t(x))(W_t - \mu_t(x'))|W_{t-1} = \sigma_{t-1}] = \text{diag}[\sigma^2 \sigma_{x}(x(i) - x'(i))]_{i=1,2},$$

where $\text{diag}[\phi(i)]_{i=1,2}$ denotes the $2 \times 2$ diagonal matrix with $\phi(i)$ being the $i$th element on the diagonal, $\sigma$ is a constant and $\sigma_{x}(\epsilon) = -0.006 + 1.006 \exp(-\epsilon/337000)$. The authors in [15, 16] also provide a recursive method to compute $s_\Sigma(w_{t-1})$ and $\Sigma_t(x|s_\Sigma(w_{t-1}))$ for this model.

3. BAYESIAN OPTIMAL DESIGN

Bayesian statistical methodology relies on the fact that all variables of the problem, such as state, parameters or observations if any, are modeled as random variables. Moreover, prior distributions are assigned to each one of them, which encapsulates any prior information or expert knowledge not captured by the model. Then inference is performed using the joint or marginal posterior appropriately. As such information is not likely to be easily available, vague prior distributions can be used. In this section we show how SA can be used to maximize the Bayesian interpretation $J_{x0}(\theta)$. Applying these ideas in the area of Air Traffic Management has been initiated by...
3.1. Bayesian inference on times of arrival

More specifically, in the case of the stopping time $T$, which is already defined as a random variable, assigning a prior $p(\tau)$ can emphasize the search or sampling in more practically useful regions, for example by considering only cases where $T < \bar{\tau}$. In this case, the simplest choice of the prior $p(\tau)$ can be the uniform distribution $\mathcal{U}_{[1, \bar{\tau}]}$. In the remainder of this paper we shall use this prior, as typically one might be able to deduce extreme or worst-case time of arrivals $\bar{\tau}$. Of course this is not the only possible choice. Other possibly diffuse distributions with most of its density mass concentrated in $[1, \bar{\tau}]$ can be considered, e.g. a negative binomial, without actually requiring $T < \bar{\tau}, \mathbb{Q}_{x_0, 0}$-almost surely (a.s.). From a Bayesian perspective we would be actually interested in actually sampling $\tau$ from the following posterior:

$$\tilde{\pi}_{T, \theta}(\tau) \propto \pi_{T, \theta}(\tau) p(\tau).$$

As expected, for the specific choice of prior being a uniform distribution, sampling from $\tilde{\pi}_{T, \theta}(\tau)$ will not return any samples $\tau > \bar{\tau}$.

In the context of Bayesian optimal design for the parameter $\theta$ we would be interested instead in minimizing the a posteriori mean of the time of arrival with respect to $\tilde{\pi}_{T, \theta}(\tau)$ rather than (18), thus yielding an alternative a posteriori reward criterion to be maximized given by

$$\tilde{J}_{x_0}(\theta) = \int u(\tau) \tilde{\pi}_{T, \theta}(\tau) \, d\tau,$$  \hspace{1cm} (18)

It is clear that we are modifying the problem formulation of Section 2.2 here. This does not involve the particular choice of prior, which implies $T < \bar{\tau}, \mathbb{Q}_{x_0, 0}$-a.s.; if a negative binomial prior was used instead this would not be the issue, but the shape of the posterior would be a result of twisting the shape of $\pi_{T, \theta}(\tau)$ by multiplying it with the prior distribution. This modification of the problem to be solved is well justified in the Bayesian statistics literature (see [1] for a few examples) and in some sense we use the prior as a flexible means of inserting valuable information not captured by modeling the times of arrival strictly as the hitting time of some set $B^i$ by the marginal of a Markov process $x_t^i$.

In Algorithm 1 we present a sequential IS algorithm that samples from $\tilde{\pi}_{T, \theta}(\tau)$ for a fixed $\theta$. First we generate samples $\tau^i \sim \tilde{\pi}_{T, \theta}$ for every $i$ and then we set the sample of the stopping time to be $\tau = \max_i \tau^i$. A convenient by-product of Algorithm 1 is that we can reuse the samples from $\tilde{\pi}_{T, \theta}(\tau)$ to obtain estimates of the probability of safe separation $\text{PSS}_\tau$. Noting that $\eta_k(G_k) \approx 1/\lambda_k$ and using equation (12) we obtain the estimate $\text{PSS}_\tau = \prod_{k=1}^{L} (1/\lambda_k)$.

To improve on the variance of this estimator one may use $L >> 1$ iid samples of Algorithm 1 and use the Monte Carlo estimator $\hat{\text{PSS}}_\tau = 1/L \sum_{l=1}^{L} \prod_{k=1}^{l} (1/\lambda_k^{(i)})$. When $\text{PSS}_\tau \geq 1 - \varepsilon$, then $\theta$ can be regarded as a feasible parameter.

Clearly, the specific choice of parameterization for the policy is crucial to obtain samples efficiently, without wasting too much sampling effort in regions of the state space where the probability of aircraft $i$ reaching $B^i$ is extremely low.

3.2. Simulated annealing

As far as $\theta$ is concerned, so far it has been treated as a deterministic parameter of the policy that needs to be computed. Using a Bayesian approach, we will treat it as a random variable and assign it a prior with density $p(\theta)$. In our application a possible choice of $p(\theta)$ can be a density peaked where each $h^1, \theta_k^0$ are near $\phi_0^i$ and zero, respectively.
Algorithm 1 Obtaining samples $\tau^i \sim \pi_T, \theta^i$ and $\tau \sim \pi_T, \theta$ given $\theta, x_0, \tau$; $i = 1, \ldots, M$.

1: Initialize permitted number of restarts at each $k = 1, \ldots, \tau$: set all $x_k = 1$, $x_{\text{max}} = \tau$.
2: for $k = 1, \ldots, \tau$ do
3: Sample $x_k \sim f_\theta (|x_{k-1})$
4: if $I_{x_k} \in \mathcal{S}$ then
5: Initialize $\mathcal{S} = 1, \ldots, M$.
6: for each $i \in \mathcal{S}$ do
7: if $I_{x_k} \in (B') = 1$ then
8: set $\tau^i = k$, augment $x^i = [x^i_{k-1}, x^i_k]$, set $\mathcal{S} \leftarrow \mathcal{S} \setminus \{i\}$ and return $\tau^i, x^i$
9: else
10: augment $x^i_k = [x^i_{k-1}, x^i_k]$ and proceed to $k + 1$
11: end if
12: end for
13: else if $k \neq \tau$ and $x_k \leq \tau$ then
14: set $x_k = x_k + 1$, go to step 3
15: end if
16: terminate without sample.
17: end if
18: end for
19: Set $\tau = \max_i \tau^i$ and return $\tau, \hat{PSS}_\tau = \prod_{k=1}^\tau \frac{1}{x^i_k}$. 

In the standard Bayesian interpretation of SA, we are interested in sampling $\theta^i$ from a density $\pi_\gamma$

$$\pi_\gamma (\theta) \propto p(\theta) f_\gamma \dot{J}_{x_0} (\theta) \gamma,$$

where $p(\theta)$ is the pdf of an arbitrary prior distribution, which must contain the maximizers $\theta^*_n$. Under weak assumptions, as $\gamma \to \infty$, $\pi_\gamma (\theta)$ becomes concentrated on the set of maximizers of $J$ [21].

Alternatively, we can introduce $\gamma_n$ artificial replicates of $\tau$, all stacked into a joint variable $\tau_1: \gamma_n$ and define the following joint posterior density:

$$\pi_n (\theta, \tau_1: \gamma_n) \propto p(\theta) \prod_{j=1}^{\gamma_n} u(\tau_j) \pi_T, \theta (\tau_j)p(\tau_j) \quad (19)$$

We will denote the marginal posterior as $\pi_n (\theta)$, for which it is easy to show that it is indeed proportional to $\pi_\gamma$, i.e.

$$\pi_\gamma (\theta) \propto \int \pi_n (\theta, \tau_1: \gamma_n) d\tau_1: \gamma_n.$$

$\{\gamma_n\}_{n \geq 0}$ is a strictly increasing integer infinite sequence, which will play the role of the inverse temperature. Under weak assumptions, for a logarithmic schedule and large $n$, one can obtain formal convergence results, and as $\gamma_n \to \infty$, $\pi_n$ becomes concentrated on the set of maximizers of $J_{x_0} (\theta)$, [8]. It has been shown empirically that logarithmic schedules can lead to slow convergence and therefore more quickly increasing rates, e.g. linear, are used. In addition, as infinite sequences are impossible to use, finite ones, $\{\gamma_n\}_{n=1, \ldots, n_{\text{max}}}$, are used instead. Also note that iteration index $n$ should not be confused with time index $t$. 

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DOI: 10.1002/acs
4. SIMULATION-BASED OPTIMIZATION

In general it is impossible to sample directly from \( \pi_n \), hence various Monte Carlo simulation schemes have been proposed. In [5, 8] this is achieved by MCMC. Instead, in this paper we propose to adopt the Sequential Monte Carlo (SMC) samplers framework of [22]. Our approach is almost identical to the application of SMC samplers presented in [6] for Maximum Likelihood parameter estimation and later in [7] for Bayesian decision theory. Later in this section we shall propose an algorithm similar to [6] and [7], which differs only in the application and context used.

4.1. SMC samplers

SMC methods are a sequential implementation of importance sampling to sample from a general distribution with density \( \pi(y) = \nu(y)/Z \), where \( \nu(y) \) can be evaluated point-wise and the normalization constant \( Z \) is unknown. A swarm of samples \( \{y^{(l)}_n\}_{l=1}^L \), called particles, evolves toward the distribution of interest by passing through a sequence of intermediate bridging distributions of increasing dimensions with densities \( \{\pi_n(y_{1:n}) = \nu_n(y_{1:n})/Z_n\}_{n \leq T} \), such that \( y = y_{0:T} \) and \( \pi_T = \pi \). This is implemented through a procedure of sampling via importance sampling and resampling (SISR), which provide appropriate approximations \( \{\hat{\pi}_n(y_{0:n})\}_{n \leq T} \). A detailed exposition of the background theory can be found in [2]. In the context state estimation or optimal Bayesian filtering for hidden Markov models SMC algorithms have also been known as Particle Filters [3]. In the most general setting SMC extends beyond this popular context [4, 22].

SMC samplers [22] are a generalization of traditional SMC methods, such as particle filtering [3]. Given a sequence of distributions with \( \{\pi_n\}_{n \geq 0} \) defined on the same or different spaces, the aim is to generate samples \( \{y^{(l)}_n\}_{l=1}^L \) from \( \pi_n \), constructed based on samples from \( \pi_{n-1} \). The core idea is to use an artificial auxiliary distribution on a space of increasing dimensions, whose density \( \{\bar{\pi}_n\} \) admits the density of interest \( \pi_n \) as its marginal. More specifically let,

\[
\bar{\pi}_n(y_{1:n}) = \pi_n(y_n) \prod_{k=n-1}^{1} \bar{q}_k(y_k | y_{k+1}),
\]

where \( \bar{q}_k \) is an appropriate backward Markov density. Then traditional SMC methods, e.g. SISR, can be used to sample from \( \bar{\pi}_n(y_{1:n}) \). New samples are proposed according to a forward Markov transition density \( q_n \). These particles are weighted recursively as follows

\[
\frac{w_n(y_{n-1:n})}{w_{n-1}(y_{n-2:n-1})} \propto \frac{\pi_n(y_n)\bar{q}_{n-1}(y_{n-1} | y_n)}{\bar{\pi}_{n-1}(y_{n-1})q_n(y_n | y_{n-1})}
\]

to provide approximations

\[
\hat{\pi}_n(dy) = \sum_{l=1}^{L} W^{(l)}_n \delta_{y^{(l)}_n}(dy) \quad \text{with} \quad W^{(l)}_n = \frac{w_n(y^{(l)}_n)}{\sum_{l=1}^{L} w_n(y^{(l)}_n)}. \]

Then a resampling procedure is introduced to copy (or multiply) particles with high weights and therefore discarding particles with low weights. The resampling procedure serves to focus the computational effort on the ‘promising’ regions of the state space. The simplest resampling scheme is multinomial resampling: generate \( L \) independent samples from \( \hat{\pi}_n(y_n) \). We shall be referring to \( \{y^{(l)}_n, W^{(l)}_n\}_{l=1}^L \) as the particle approximation \( \hat{\pi}_n \) to \( \pi_n \).
Many sharp convergence results are available for SMC algorithms; under fairly weak assumptions one can obtain almost sure convergence results with respect to the population size \( L \)

\[
\sum_{l=1}^{L} W_n^{(l)} \delta_{y_n^{(l)}}(dy) \overset{L \to \infty}{\overset{a.s.}{\longrightarrow}} \pi_n(dy),
\]

as well as central limit theorems, see [22] for SMC samplers and [2] for a book length review. For more details on traditional SMC and SMC samplers we refer to [3, 4, 22], respectively.

4.2. SMC samplers for optimal policy design

In our case, we set \( y_{1:n} = (\theta, \tau_{1:\gamma_n}) \) and \( y_n = (\theta, \tau_n) \) with \( \tau_n = \tau_{\gamma_n+1:\gamma_n} \). A possible candidate for \( \pi_n \) can take the form of (19), but we shall further generalize that formulation to include more flexible non-integer annealing schedules following the ideas in [6, 7]. Instead of (19) we shall use

\[
\pi_n(\theta, \tau_{1:\gamma_n}) = p(\theta) \prod_{j=1}^{\gamma_n} u(\tau_j) \pi_{T, N}(\tau_j) p(\tau_j)^{v_n}, \tag{20}
\]

where \( \{\gamma_n\}_{n=1,\ldots,n_{\text{max}}} \) is a monotone increasing integer sequence with \( \gamma_n \in \mathbb{N}^+ \), and \( \{v_n\}_{n=1,\ldots,n_{\text{max}}} \) is a real sequence such that \( v_n \in [0, 1] \), with the only restriction being that \( v_{n_{\text{max}}} = 0 \). Note that when \( v_n = 0 \), the marginal \( \pi_n(\theta) \) remains as before.

To obtain samples for this specific distribution sequence \( \{\pi_n(\theta_n, \tau_n)\} \), we propose to propagate \( \theta \) at each \( n \) using a random walk transition density for \( q_n(\theta_n|\theta_{n-1}) \) and then once \( \theta_n \) is sampled, sample the hitting times from the true dynamics, \( \tau_n \sim \prod_{j=\gamma_{n-1}+1}^{\gamma_n} \pi_{T, N}(\tau_j) \). We will also use the backward transition density

\[
\bar{q}_{n-1}(\theta_{n-1}, \tau_{n-1} | \theta_n, \tau_n) \propto q_n(\theta_{n-1} | \theta_n) u(\tau_{\gamma_{n-1}})^{v_{n-1}},
\]

which serves as an approximation that can resemble closely the optimal backward transition density with respect to the variance of the weights presented in [22]. The weights can then be propagated according to

\[
\frac{w_n(\tau_n)}{w_{n-1}(\tau_{n-1})} \propto u(\tau_{\gamma_n})^{v_n} \prod_{j=\gamma_{n-1}+1}^{\gamma_n} u(\tau_j).
\]

In Algorithm 2 we summarize the proposed SMC algorithm which can be used to maximize \( J_{x_0}(\theta) \) for each aircraft.

4.3. Implementation notes and related approaches

The resampling step adds some variance to the weights to keep the particle set alive, i.e. to prevent all particles having negligible weight apart very few or even a single one. In addition, a key element for the success of any proposed SMC algorithm is that the shape of \( \pi_n \) should be close to that of \( \pi_{n-1} \). This is reflected in the variance of the importance weights. Therefore, it is important to carefully choose \( \bar{q}_n \) and \( q_n \) to achieve good mixing properties and keep the variance of the importance weights small. As the resampling step tends to contribute to an increase in the variance of the weights it has been proposed to use resampling only when it is necessary, i.e. when the variance
Algorithm 2 Simulation-based trajectory planning.

1: Initialization, at \( n = 0 \): for each particle \( l = 1, \ldots, L \): sample from the prior \( \theta_{0}^{(l)} \sim p(\cdot) \) and assume \( W_{l}^{1} = \frac{1}{L} \).
2: for \( n = 0, \ldots, n_{\text{max}} \): do
3:   for each particle \( l = 1, \ldots, L \): do
4:     if \( n \geq 1 \), sample control parameter particles \( \theta_{n}^{(l)} \sim q_{n}(\cdot|\theta_{n-1}^{(l)}) \)
5:     for \( j = \gamma_{n-1} + 1, \ldots, \gamma_{n} \) using \( \theta_{n}^{(l)} \) and Algorithm 1, sample \( \tau_{j}^{(l)} \)
6:     Weighting particles: set \( w_{n}^{(l)} = W_{n}^{(l)} u(\tau_{j}^{(l)})^{\gamma_{n}} \prod_{j=\gamma_{n-1}+1}^{\gamma_{n}} u(\tau_{j}^{(l)}) \).
7:     if Algorithm 1 returns no sample, set \( w_{n}^{(l)} = 0 \).
8:     end for
9:   Normalize weights: \( W_{n}^{(l)} = w_{n}^{(l)} / \sum_{l'} w_{n}^{(l')} \).
10:  Resample, if \( \text{ESS}_{n} \leq \text{L}_{\text{thresh}} \), to get new particle set \( \{(\theta_{n}^{(l)}, \tau_{n}^{(l)})\}_{l=1}^{L} \) with equal weights \( W_{n}^{(l)} = \frac{1}{L} \).
11: end for
12: Compute the estimate of optimal \( \theta \)

of the unnormalized weights exceeds some threshold. This can be implemented by monitoring the effective sample size (ESS) [4]

\[
\text{ESS}_{n} = \left( \frac{1}{L} \right) \sum_{l=1}^{L} (W_{n}^{l})^{2}
\]

and resampling only if it drops under some pre-specified value, e.g \( \frac{1}{2} \). The ESS is always less than \( L \) and can be interpreted as the number of perfect samples from the target distribution that would yield the same estimator variance. Note that we expect that Algorithm 2 will yield lower ESS compared with the ones found in [6, 7] or in other more traditional applications of SMC such as nonlinear filtering. This is because the algorithm will be rejecting samples which violate constraints, although hopefully the frequency of the latter should reduce with \( n \).

A common criticism can be that as in any annealing algorithm, some tuning might be required for the annealing schedule. The tradeoff is clearly between fast convergence and the chance to be trapped in local optima. The non-integer annealing schedules \( \{\gamma_{n}, \nu_{n}\}_{n=1}^{n_{\text{max}}} \) can make tuning slightly easier, and it is also possible to adapt \( \gamma_{n}, \nu_{n} \) to the drop in the ESS and use a slower annealing schedule to spend more time exploring areas where the algorithm is prone to get stuck. Similar ideas have been used recently in the context of approximate Bayesian computation in [23].

Finally, as we are interested in estimating the maximizers \( \theta^{*} \), using the mean of \( \hat{\theta}_{n_{\text{max}}} \) can be a bad estimate in the case of multi-modalities. Therefore, one should be careful to use the mode of \( \hat{\theta}_{n_{\text{max}}} \) or the mean of some cluster of \( \{\theta_{n_{\text{max}}}^{l}\}_{l=1}^{L} \) with significant weight as estimates \( \theta^{*} \).
our approach is the explicit handling of constraints, which has not appeared earlier in [6, 7]. Moreover, readers familiar with the stochastic programming literature [24] might find some resemblance to our approach. Nevertheless, we clarify that our approach is completely different as we are approximating the solution of the true problem, i.e. estimating the true maximizer, rather than solving an approximate problem based on sample averages. In addition we avoid the restrictive setting of iid sampling and convex problems very common in the literature of stochastic programming.

5. SIMULATIONS

We begin by discussing the details of the wind field model discussed in Section 2.3. We use RUC2 prediction gridded data for flight level (FL) 300 over the Continental United States (CONUS). The data was a wind velocity forecast for time 20:00, 21 of February 2008, and was downloaded from the NOAA server three hours earlier. In order to perform interpolation with respect to time we used the same forecast for time 23:00. Each grid consists of 451 longitudinal and 337 points, where each point is separated by 13.545 km from others either horizontally or vertically. In Figure 1 we plot a velocity diagram of the wind velocity for time 20:00 and in Figure 2 we plot the $U$- and $V$-components separately with respect to the longitude and latitude.

We will construct an artificial conflict involving 4 aircraft heading toward each other in the rectangle defined by 255–260° in longitude and 30–35° in latitude and refer to this area as the area of conflict. The specific part of CONUS was chosen due to the severity of the winds appearing in that area during the time interval considered, which results in a harder conflict resolution problem. Aircraft are initialized in the different corners of the rectangle and are heading diagonally toward the opposite corner. We will assume that all aircraft are cruising at the same constant speed, 253 ms$^{-1}$, which is a typical value for the specific flight level [11]. In the remainder of this section, we shall present the performance of Algorithm 2 for each of the presented dual mode policies in Section 2.1.2, i.e. when fixed heading or tracking is used for the last block. Algorithm 2 will be implemented using a linear annealing

![Figure 1. Velocity diagram of wind velocity of a snapshot of the data set recorded at time 20:00, 21 of February 2008: (a) velocity diagram of wind over CONUS. The rectangle is the area where a conflict appears, which we aim to resolve. A magnified plot of this area is plotted in Figure 1(b) and (b) $U$- and $V$-components of wind velocity over CONUS at time 20:00, 21 of February 2008, with respect to latitude and longitude.](image-url)
Figure 2. U- and V-components of wind velocity over CONUS at time 20:00, 21 of February 2008. with respect to latitude and longitude.

schedule, with \( n_{\text{max}} = 100, L = 1000, q_n \) a Gaussian kernel with decaying bandwidth, \( \sigma^2 = 5 \) and for the policy we set \( H = 4 \) and \( p = 3 \).

5.1. Fixed heading for mode 2

Here we consider the case where in the final block the heading consists of the angle to exit after the headings in block \( H - 1 \) have been applied. As explained in Section 2.1.2, the control inputs in the first \( H - 1 \) blocks depend on \( \theta \), which is computed using Algorithm 2. The second mode of the policy is block \( H \), which is independent of \( \theta \). Nevertheless, it affects the performance of the policy significantly in terms of the closeness to optimality, the statistical efficiency of Algorithms 1 and 2, as well as the sensitivity to constraints and size of target sets. Also the second mode is crucial for the resulting estimates of \( \theta^* \) and the underlying state trajectories, which in turn inherit some characteristic features due to the specific nature of the second mode.

In Figure 3(a) we illustrate the resulting average aircraft trajectories after Algorithm 2 has been executed. One can see that there is a mild form of symmetry in the direction each aircraft turns to. We stress that turning in the same direction is not imposed by the policy or the algorithm, but is a result of the simulation-based selection procedure in the algorithm. In Figures 3(b) and (c) we plot the approximate cost and the ESS with respect to iteration \( n \). Finally in Figure 3(d) we plot the minimum, maximum and mean separation of each pair of aircraft, which are computed by using the estimated optimal policy in Algorithm 1 to obtain independent Monte Carlo samples of the aircraft trajectories.

It is clear that when fixed heading is used for the second mode of the policy, separation constraints are easier to handle, at the cost of requiring larger exit zones \( B^i \) to achieve reasonable performance and maintain reasonable sampling efficiency. This is because after the first mode of the policy completes, the aircraft turns toward the exit point, but then drifts away from it due to the unexpected wind.

5.2. Tracking for mode 2

On the other hand when tracking is used, although the exit zones can be made arbitrarily small, the resulting trajectories will be closer to each other and the policy seems to be more sensitive to constraints. This was confirmed by our initial simulations. So that a fair comparison can be made when Algorithm 2 is used, we used a gradual smooth constraint-tightening approach. We used a varying minimum separation for each iteration \( n \), which was given by \( d_n = d/(0.8 + n^{-0.5}) \). In fact, \( d_{n_{\text{max}}} < d \), but the values are quite close. This allows comparisons to be
made on a fair level, as the target zones for each aircraft used in our simulations were better suited for the fixed heading case.

In Figure 4 we show the same plots as in Figure 3 for the same problem, when tracking is used in the second mode of the policy together with gradual tightening of the constraints. We observe that the approximate expected costs this time are marginally lower but similar compared to the fixed heading case. As regards the sampling efficiency monitored by the ESS, both policies seem to perform similarly.
Figure 4. Performance of Algorithm 2 when tracking is used in mode 2: for more details we refer to the caption of Figure 3. Simulations lasted 89 min when implemented in Matlab: (a) trajectories of resulting maneuvers; (b) approximate expected cost in minutes with iteration \( n \); (c) effective sample size with iteration \( n \); and (d) separation for every pair of aircraft with time.

6. CONCLUSIONS

In this paper we considered the ATM problem of conflict resolution for level flight, when constant speed is used and only changing the heading of the aircraft is allowed at each time. We proposed a conveniently parameterized dual-mode policy to automatically design maneuvers using simulation, so that each aircraft reaches a target in the expected minimum time of arrival in the presence of stochastic wind disturbances modeled by a realistic model derived from the actual forecast data. We showed how effective our approach is by appropriate simulations. As regards the resolution properties, an implicit feature of the presented policy is that each aircraft is forced earlier rather than later to act preemptively by executing maneuvers in case a conflict arises. If one takes into account that the actual execution times of our implementations can be drastically reduced by at least one order of magnitude using parallel implementation GPUs [10], submitted), this makes our proposed approach suitable for mid/short term conflict resolution with a horizon of the order of tens of minutes.
ACKNOWLEDGEMENTS

This work is supported by the European Commission under project iFly and EPSRC under grant numbers FP6-TREN-037180, EP/H021558/1. The authors are grateful to the anonymous reviewers for their suggestions to improve the quality and presentation of this paper.

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