# Algorithms for the incorporation of predictive information in surveillance theory†

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The paper discusses the problem of tracking a platform given probabilistic information concerning its possible destinations. Using a bayesian approach, we derive an algorithm which optimally incorporates the real-time measurements of the platform position with the predictive information about its destination. This algorithm can be efficiently implemented using some results from optimal smoothing. The results are illustrated by a simple example.

#### 1. Introduction

An important problem in the current U.S. Navy system is the design of an overall surveillance system. This surveillance system has two primary objectives the first is to keep track of all the ships located in a given area; the second is to maintain a record of all the information related to these ships (status, nationality, possible destination, etc.). The information which is available to such a surveillance system consists of a combination of the following data:

- (i) the identification of the ship's class and its dynamics;
- (ii) some information about the ship's origin, its destination and its refuelling stops;
- (iii) sightings by other ships (submarines or surface ships), by aeroplanes or by satellites; and
- (iv) some sonar information obtained from tracking sonar arrays.

A close look at the available data indicates that it can be divided into two categories. Sightings by ships, aeroplanes, sonar or radar, and the ship's origins provide information about the past of the ship's trajectory, whereas information about the ship's destination and refuelling stops provide information about the future of the ship's trajectory. In this paper, we derive an algorithm which can be used to incorporate both categories of information, in order to obtain an overall best estimate of the ship's position.

Our basic approach is to assume that the ship's dynamics are governed by a state-

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space model of a stochastic dynamical system, of the form

$$dx(t) = A(t)x(t) dt + dU(t)$$
(1)

where U(t) is a Wiener process with covariance

$$E\{U(t) U'(s)\} = \int_{0}^{\min(t,s)} Q(\tau) d\tau \tag{2}$$

Models of this form have been proposed by a number of authors (Singer 1970, Singer and Behnke 1971, Morgan 1976) for the modelling of ship trajectories with and without autopilots, or in manoeuvring situations. The detailed modelling of ship dynamics is not important to the subsequent results; the reader is referred to the above for further discussions of these issues.

The model in (1) describes the ship dynamics before any additional information on its destination is known. In the next section, we describe an algorithm which modifies (1) to obtain a new markovian model which incorporates the predictive information concerning a single destination.

This algorithm uses the basic approach of obtaining backwards markovian models, as given in Verghese and Kailath (1979). Based on this new model, the optimal estimate of ship position using past measurements can be obtained from standard estimation algorithms such as Kalman filtering (Van Trees 1968).

In § 3, we extend the results of § 2 to incorporate probabilistic predictive information concerning multiple possible destinations. The optimal estimate of the ship's position, based on past and predictive information, is computed using a multiple-model approach (Athans et al. 1977). Furthermore, a likelihood-ratio method to identify the validity of the intelligence information is presented, leading to an efficient implementation of the tracking algorithm using fixed-interval smoothing formulae (Mayne 1966, Fraser 1967, Wall et al. 1981). In § 4, we discuss how this method can be adapted for manoeuvre detection using the results of Willsky and Jones (1978). Section 5 contains an example which illustrates the basic algorithm.

# 2. Incorporation of predictive information

Assume that a state-space model of a ship's trajectory is given by eqns. (1)–(2) and the initial condition

$$x(0) \sim N(m_0, \Pi_0) \tag{3}$$

where x(0) is independent of U(t) for all t in [0, T]. Assume that we receive the additional predictive information that, at time T, the ship's position will be observed probabilistically, in terms of a mean value and an error, as

$$x_T \sim N(m_T, \, \Pi_T) \tag{4}$$

Our goal is to obtain a new model which already incorporates this predictive information. Mathematically, we are interested in obtaining a new description of the trajectories of the process, of the form

$$dx(t) = \bar{A}(t)x(t) dt + d\bar{U}(t)$$
(5)

$$x(0) = N(\bar{x}_0, \bar{\Sigma}_0) \tag{6}$$

such that x(0) is independent of  $\overline{U}(t)$  for all t in [0, T], and such that the probability distribution induced by (5) and (6) on the space of sample paths  $C\{[0, T]; \mathbb{R}^n\}$  is equal

to the probability induced by (1), (2) and (3), conditioned on knowledge of the predictive information (4). The first step in this process is to describe the predictive information as an observation. Let

$$m_T = x_T + v_T \tag{7}$$

where  $v_T \sim N(0, \Pi_T)$  is independent of U(t) for t in [0, T]. Let  $(\Omega, F, P)$  be our basic probability space, where  $x_0, v_T$  and  $U(s), 0 \le s \le T$  are defined. Define the two increasing sequences of  $\sigma$ -fields

$$G_{t} = U_{t} \vee x_{0} F_{t} = U_{t} \vee X_{0} \vee M_{T}$$
,  $0 \le t \le T$ 

where  $U_t = \sigma\{U(s), \ 0 \le s \le t\}$  is the sequence of  $\sigma$ -fields generated by the Wiener process  $U(\cdot)$ , and where  $X_0$  and  $M_T$  are the  $\sigma$ -fields generated by x(0) and  $m_T$ . Since the noise  $U(\cdot)$  is independent of x(0), the process  $(U(t), G_t)$  is a martingale; that is

$$E\{U(s) - U(t)|G_t\} = 0$$

for  $0 \le t \le s \le T$ . Hence, we can write (1), in integrated form, as

$$x(t) = x(0) + \int_0^t Ax(s) \, ds + \int_0^t dU(s)$$
 (8)

where x(t) is a Markov process, the middle term is  $G_t$ -predictable, and the last term is a  $G_t$ -martingale. A model satisfying this property is said to be a markovian model of  $x(\cdot)$ . Most filtering or estimation results have been obtained for processes described by markovian models.

Note that the process U(t) is not a martingale with respect to  $F_t$  (indeed, U(t) and  $m_T$  are correlated). Therefore, the problem of incorporating the information (4) in the model (1) and (2) is essentially the one of constructing a markovian model of  $x(\cdot)$  with respect to  $F_t$ , i.e. with respect to all the a priori information available, including predictive information.

To do so, we shall follow a method similar to the one used by Verghese and Kailath (1979) for the derivation of backwards markovian models. The first step is to construct a Doob-Meyer decomposition of U(t) with respect to  $F_t$ , i.e.

$$U(t) = \int_0^t E\{dU(\tau)|F_\tau\} + \tilde{U}(t)$$
 (9)

where

$$E\{\tilde{U}(s) - \tilde{U}(t)|F_t\} = 0 \text{ for } 0 \le t \le s \le T$$

The reader is referred to Meyer (1966) and Wong (1973) for a description of the properties of such a decomposition. From (9) it is clear that to obtain the martingale  $\tilde{U}(t)$ , one need only compute the predictable projection  $E\{dU(t)|F_t\}$ , where

$$dU(t) \triangleq U(t+dt) - U(t)$$

However, before doing so, it will be useful to decompose  $F_t$  as

$$F_t = U_t \vee X_0 \vee \tilde{X}(t, T) \tag{10}$$

where  $\phi(\cdot, \cdot)$  denotes the transition matrix of  $A(\cdot)$ , so that

$$\tilde{x}(t, T) = x(t) - \phi(t, T)m_T = \int_t^T \phi(t, s) \ dU(s) - \phi(t, T)v_T$$
 (11)

Let  $\tilde{X}(t, T)$  denote the  $\sigma$ -field generated by  $\tilde{x}(t, T)$ . This  $\sigma$ -field, by construction, is independent from  $G_t = U_t \vee X_0$ . To justify the decomposition (10) note that

$$x(t) \in U_t \vee X_0 \tag{12}$$

and also observe that

$$x(t) \vee \tilde{x}(t, T) = X(t) \vee M_T \tag{13}$$

(i.e. the knowledge of  $(x(t), \tilde{x}(t, T))$  is equivalent to the knowledge of  $(x(t), m_T)$ ). We are now in a position to prove our main theorem.

Theorem 1: Modified model

Let  $\Pi(t, T)$  denote the solution of the equation

$$-\frac{d}{dt}\Pi(t, T) = -A\Pi - \Pi A' + Q$$

$$\Pi(T, T) = \Pi_{T}$$
(14)

Then, the model

$$x(t) = x(0) + \int_{0}^{t} Ax(s) ds + \int_{0}^{t} Q(s)\Pi^{-1}(s, T)[\phi(s, T)m_{T} - x(s)] ds + \tilde{U}(t)$$
(15)

is a markovian model of  $x(\cdot)$  with respect to  $F_t$ , where  $(\tilde{U}(t), F_t)$  is a brownian motion with intensity  $\int_0^t Q(s) ds$ .

Proof

The independence of G, and  $\tilde{x}(t, T)$  gives

$$E\{dU(t)|F_t\} = E\{dU(t)|G_t\} + E\{dU(t)|\tilde{X}(t,T)\} = E\{dU(t)|\tilde{x}(t,T)\}$$

because  $U_t$  is a  $G_t$  martingale and  $\tilde{X}(t, T)$  is generated by  $\tilde{x}(t, T)$ . Now, dU(t) and  $\tilde{x}(t, T)$  are gaussian random variables, so that

$$E\{dU(t)|\tilde{x}(t,T)\} = E\{dU(t)\tilde{x}'(t,T)\}E\{\tilde{x}(t,T)\tilde{x}'(t,T)\}^{-1}\tilde{x}(t,T)$$
(16)

From (11), we get

$$E\{dU(t) \ x'(t, T)\} = -Q(t) \ dt$$

$$E\{\tilde{x}(t, T) \ \tilde{x}'(t, T)\} = \int_{t}^{T} \phi(t, s) Q(s) \phi'(t, s) \ ds + \phi(t, T) \Pi_{T} \phi'(t, T)$$
(17)

Note that (17) is the solution of (14), so that

$$E\{\tilde{x}(t, T) \, \tilde{x}'(t, T)\} = \Pi(t, T)$$

Hence

$$E\{dU(t)|F_t\} = Q(t)\Pi^{-1}(t, T)(\phi(t, T)m_T - x(t)) dt$$
(18)

Equation (18) represents the predictable part of the Doob-Meyer decomposition of U(t) with respect to the  $\sigma$ -fields  $F_t$ . Since U(t) has continuous sample paths, the quadratic variation of  $\tilde{U}(t)$  is the same as the quadratic variation of U(t), so it is an  $F_t$  Wiener process with covariance

$$E\{\widetilde{U}(t)\ \widetilde{U}(s)\} = \int_{0}^{\min(t,s)} Q(\tau)\ d\tau$$

Now, substituting (18), into the integrated form of (1) yields the results of the theorem.

Remark

The model (15) is the same as the backwards markovian model obtained by Verghese and Kailath (1979, see also Ljung and Kailath (1976) and Lainiotis (1976)). However, the model (15) is a *forwards* model. The martingale method that we have used here to construct (15) is very general and applies also to a large class of problems in the study of random fields (Willsky 1979).

The results of Theorem 1 provide us with a model which describes the evolution of the sample paths of x(t) from an initial condition x(0). In order to properly specify the distribution of sample paths, we must provide a description of the initial distribution for x(0). This initial distribution will be the conditional probability of x(0), given the prior statistics

$$x(0) \sim N(m_0, \Pi_0)$$

and the posterior observation

$$m_T = x_T + v_T$$

From (11), we can transform this observation to

$$\phi(0, T)m_T = x(0) - \tilde{x}(0, T) \tag{19}$$

where  $\tilde{x}(0, T)$  is independent of x(0), with distribution

$$\tilde{x}(0, T) \sim N(0, \Pi(0, T))$$

Hence, the a posteriori distribution of x(0), given  $M_T$ , is

$$x(0) \sim N(\hat{x}_{\epsilon}(0), \Pi_{\epsilon}(0))$$
 (20)

where

$$\Pi_{s}(0) = (\Pi_{0}^{-1} + \Pi^{-1}(0, T))^{-1}$$
(21)

$$\hat{x}_s(0) = \Pi_s(0)(\Pi_0^{-1}m_0 + \Pi^{-1}(0, T)\phi(0, T)m_T)$$
(22)

This last expressiom corresponds to the two-filter Mayne-Fraser smoothing formula (Mayne 1966, Fraser 1967) obtained from the observations (3) and (4).

Based on these initial conditions, the correspondence between the model in (1) and (2), with initial condition (3) and predictive observation (4), and the model in Theorem 1 can be stated precisely.

Theorem 2

Consider the markovian model of Theorem 1, with initial distribution (20). The probability measure defined on the sample paths  $x(\cdot)$  by this model is equal to the probability measure defined on the sample paths  $x(\cdot)$  by the model (1) and (2) with initial condition (3), conditioned on the  $\sigma$ -field  $M_T$ , almost surely.

Proof

Let  $P_1$  be the distribution induced by (1) and  $P_2$  the distribution induced by (15). By construction, both models yield continuous sample paths with gaussian statistics. Hence, it is sufficient to show that the mean and covariance of  $x(\cdot)$  under both models is identical. By (20), we know

$$E_{P_1}\{x(0)|M_T\} = E_{P_2}\{x(0)\} = \hat{x}_s(0) \quad \text{a.s.}$$

$$E_{P_1}\{[x(0) - \hat{x}_s(0)] [x(0) - \hat{x}_s(0)]'|M_T\} = E_{P_2}\{[x(0) - \hat{x}_s(0)] [x(0) - \hat{x}_s(0)]'\} = \Pi_s(0)$$

Furthermore

$$\begin{split} E_{P_1}\{x(t)|M_T\} &= E_{P_1}\{x(0)|M_T\} + E_{P_1}\left\{\int_0^t Ax(s) \ ds \ \bigg| \ M_T\right\} \\ &+ E_{P_1}\{U(s)|M_T\} \\ &= E_{P_1}\{x(0)|M_T\} + \int_0^t AE_{P_1}\{x(s)|M_T\} \ ds \\ &+ \int_0^t Q(s)\Pi^{-1}(s, T)[\phi(s, T)M_T - E_{P_1}\{x(s)|M_T\}] \ ds \end{split}$$

because

$$\begin{split} E_{P_1} \big\{ \tilde{U}(t) | M_T \big\} &= E_{P_1} \big\{ E_{P_1} \big\{ \tilde{U}(t) | X_0, M_T \big\} | M_T \big\} \\ &= 0 \end{split}$$

which implies that, for all t,

$$E_{P_1}\{x(t)|M_T\} = E_{P_2}\{x(t)\}$$
 a.s.

An identical argument using Ito's formula (Breiman 1968) establishes

$$E_{P_1}\{f(x(t))|M_T\} = E_{P_2}\{f(x(t))\}$$

and

$$E_{P_1}\{f(x(s))g(x(t))|M_T\} = E_{P_2}\{f(x(s))g(x(t))\}$$

for  $0 \le s$ ,  $t \le T$ , and for twice differentiable functions f and g. Hence, the finite-dimensional distributions of  $P_1$  conditioned on  $M_T$  and  $P_2$  agree, proving the theorem.

Example: The pinned case

Let  $\Pi_T = 0$ , and x(t) = m. In this case, the value of x(t) is fixed and the model (15) has the property that  $\pi(T, T) = 0$ . Thus, by denoting  $\delta(t) = x(t) - \phi(t, T)m$  one gets

$$d\delta(t) = A(t) - Q(t)\pi^{-1}(t, T)\delta(t) dt + d\tilde{U}(t)$$

where  $\pi^{-1}(t, T) \to \infty$  as t + T. This means that as  $t \to T$  the state x(t) is steered stronger and stronger towards m, and  $\delta(t) \to 0$  almost surely as  $t \to T$ .

This property is consistent with previous studies of the pinned brownian motion (cf. Breiman).

The main feature of the procedure that has been used here to obtain the model (15) with initial conditions (20)-(22) is that it is *bayesian*. This means that the model (1) and (2) has been assumed a priori for  $x(\cdot)$ . The advantage of this procedure is that it can be used to incorporate all the information available on the ship trajectory sequentially.

For example, suppose that additional predictive information concerning possible refuelling stops is available. That is, we are given

$$x(t_i) \sim N\{m_i, \ \Pi_i\} \tag{23}$$

with  $0 \le t_i \le T$  for i = 1, ..., N.

By an analogous procedure to the construction in Theorem 1, we can construct a markovian model which incorporates the above predictive information. Now, assume that real-time measurements

$$dy(t) = Cx(t) dt + dV(t)$$
(24)

are available, where the noises dU(t) and dV(t) are independent. Then, based on the markovian model of Theorem 1, the optimal estimate of x(t), given the predictive information (23) and the real-time observations (24), can be obtained from a Kalman filter based on the markovian model (15) with observation (24).

An alternative approach to obtaining a markovian model of ship trajectories which included destination information was derived in Weisinger (1978), in the context of simple brownian motion. In Weisinger (1978), a non-bayesian approach was used to incorporate the predictive information. In the case that we are considering here, a simple generalization of Weisinger (1978) would yield the model

$$d\begin{bmatrix} x(t) \\ x(T) \end{bmatrix} = \begin{bmatrix} A(t) - Q(t)\Pi^{-1}(t, T) & Q(t)\Pi^{-1}(t, T)\phi(t, T) \\ 0 & 0 \end{bmatrix} \begin{bmatrix} x(t) \\ x(T) \end{bmatrix} dt + \begin{bmatrix} d\tilde{U}(t) \\ 0 \end{bmatrix}$$
(25)

where the initial conditions are given by

$$\begin{bmatrix} x(0) \\ x(T) \end{bmatrix} \sim N \begin{pmatrix} \begin{bmatrix} m_0 \\ m_T \end{bmatrix}, \begin{bmatrix} \Pi_0 & 0 \\ 0 & \Pi_T \end{bmatrix} \end{pmatrix}$$

and where  $\tilde{U}(\cdot)$  is a Wiener process with covariance (2). The main aspect of this model is that it is obvious by replacing  $m_T$  by x(T) in (15) and by setting  $\Pi(T, T) = 0$  in the differential equation for  $\Pi(\cdot, T)$ . To see how this model differs from (15) we need only note that its transition matrix is given by

$$\psi(t, s) = \begin{bmatrix} \psi_1(t, s) & \psi_2(t, s) \\ 0 & I \end{bmatrix}$$
 (26)

where  $\psi_1(t, s) = \pi(t, T)\phi'(s, t)\pi^{-1}(s, T)$ 

$$\psi_2(t, s) = \phi(t, T) - \psi_1(t, s)\phi(s; T)$$

Then if m(t) and  $\pi(t)$  denote the mean and variance of x(t), respectively, one has

$$m(t) = \psi_1(t, 0)m_0 + \psi_2(t, 0)m_T$$

$$\pi(t) = \psi_1(t, 0)\pi_0\psi_1'(t, 0) + \psi_2(t, 0)\pi_T\psi_2'(t, 0)$$

$$+ \int_0^t \psi_1(t, s)Q(s)\psi_1'(t, s) ds$$

Thus, by noting that

$$\psi_1(0, 0) = I, \quad \psi_1(T, 0) = 0$$
  
 $\psi_2(0, 0) = 0, \quad \psi_2(T, 0) = I$ 

we see that  $x(0) \sim N(m_0, \pi_0)$  and  $x(T) \sim N(m_T, \pi_T)$  exactly. By comparison the initial conditions that were used for the bayesian model (15) were

$$x(0) \sim N(\hat{x}_s(0), \pi_s(0))$$

where  $\hat{x}_s(0)$  is a linear combination of  $m_0$  and  $\phi(0, T)m_T$ . This means that the information on the time of arrival of the ship was used to provide some information on its time of departure. Such a procedure makes sense: if a ship crosses the Atlantic and arrives on a certain date, it should leave at least one week before.

The previous comments illustrates the main advantages of bayesian methods as compared to non-bayesian methods for the fusion of information:

- (i) bayesian methods are sequential (any additional information can be incorporated in the model); and
- (ii) no piece of information is overweighted with respect to another.

# 3. Predictive information about multiple destinations

The results of Theorems 1 and 2 describe an approach for deriving markovian models which incorporate predictive information about a single possible destination. However, it is common that predictive information consists of multiple hypotheses concerning the possible destinations of the object of interest. A typical intelligence report might read: 'the ship x will be heading for point A with 75 per cent certainty; otherwise, it will head for point B'. In this section, we discuss a bayesian estimation scheme which will incorporate this information, together with some real-time measurements, to obtain the optimal state estimate.

Figure 1 illustrates the case when there exist a number of possible destinations for the ship, characterized by

$$H_i: x(t_i) \sim N(m_i, \pi_i), \quad 1 \leq i \leq N$$

For each of the possible hypotheses  $H_i$ , we can use the results of Theorems 1 and 2 to obtain the markovian models

$$dx(t) = (A_i(t)x(t) + s_i(t)) dt + \tilde{U}_i(t)$$
(27)

where

$$A_i(t) = A - Q(t)\pi_i^{-1}(t, t_i)$$
(28)

$$s_i(t) = Q(t)\pi_i^{-1}(t, t_i)\phi(t, t_i)m_i$$
 (29)

where  $\pi(t, t_i)$  is the solution of (14) with  $\pi(t_i, t_i) = \pi_i$ 

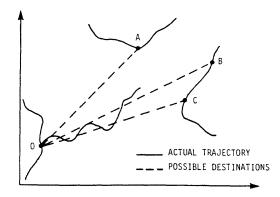


Figure 1. Determination of the ship's destination from multiple hypotheses.

The model described by (27), (28), and (29) is a markovian model of  $x(\cdot)$ , conditioned on the information  $m_i$  (that is,  $H_i$  is true), represented as

$$m_i = x(t_i) + v_i$$
$$v_i \sim N(0, \pi_i)$$

Suppose now that real-time observations are available, to discriminate the hypotheses  $H_i$ . Assume that these observations are given by

$$dy(t) = C_i(t)x(t) dt + R^{1/2}(t) dV(t)$$
(30)

where V(t) is a standard Wiener process independent of U(t). The likelihood ratio based on the sequence of measurements  $V_t = \{y(s), 0 \le s \le t\}$  is defined for each pair of hypotheses as

$$L_{ij}(t) = \frac{P(H_i|Y_t)}{P(H_i)} \frac{P(H_j)}{P(H_j|Y_t)}$$
(31)

where  $P(H_i)$  denotes the a priori probability of  $H_i$ .

Thus, evaluation of the likelihood ratio between models can be obtained from a bayesian approach by evaluating the *a posteriori* probabilities  $P_i(t) = P(H_i|Y_i)$ . These probabilities can be evaluated recursively using Kalman filters (see Dunn (1977) and Athans *et al.* (1977) for a derivation of these equations) as

$$d\hat{x}_{i}(t) = [A_{i}(t)\hat{x}_{i}(t) + s_{i}(t)] dt + \Sigma_{i}(t)C'_{1}(t)R^{-1}(t)$$

$$\cdot [dY(t) - C_{i}(t)\hat{x}_{i}(t) dt]$$
(32)

$$dP_{i}(t) = P_{i}(t) \left[ C_{i}(t)\hat{x}_{i}(t) - \sum_{j=1}^{N} P_{j}(t)C_{j}(t)\hat{x}_{j}(t) \right] R^{-1}(t)$$

$$\cdot \left[ dy(t) - \sum_{j=1}^{N} P_{j}(t)C_{j}(t)\hat{x}_{j}(t) dt \right]$$
(33)

$$\frac{d}{dt} \Sigma_i(t) = A_i(t) \Sigma_i(t) + \Sigma_i(t) A_1'(t) + Q(t)$$
$$- \Sigma_i(t) C_i'(t) R^{-1}(t) C_i(t) \Sigma_i(t)$$
(34)

where  $x_i(t)$  and  $\Sigma_i(t)$  are the conditional estimate and convariance given that  $H_i$  is true. The initial conditions are

$$\hat{x}_i(0) = \hat{x}_{si}(0)$$

$$P_i(0) = P(H_i)$$

$$\Sigma_i(0) = \pi_{si}(0)$$

where  $\hat{x}_{si}(0)$  and  $\pi_{si}(0)$  are given by (21) and (22). We can now state the main result of this section.

# Theorem 3

Denote the  $\sigma$ -field generated by the predictive information as K. Then, the conditional density of x(t) given the information K and the observations  $0 \le y(s) \le t$ , is given by

$$p(s) = \sum_{i=1}^{N} P_i(t) p_i(x, t)$$

where  $p_i(x, t)$  is a gaussian density with mean  $\hat{x}_i(t)$  and variance  $\Sigma_i(t)$ .

Proof

The proof follows directly from Bayes' formula, as, for a Boreal set A, in  $\mathbb{R}^n$ ,

$$P\{x(t) \in A \mid K, y(s), 0 \le s \le t\} = \sum_{i=1}^{N} [\text{prob} \{x(t) \in A \mid H_i, y(s), 0 \le s \le t\}$$

$$\times \text{prob} \{H_i \text{ is true} \mid K, y(s), 0 \le s \le t\}]$$

which establishes the theorem, using Theorems 1 and 2.

The eqns. (32)–(34) requires the parallel processing of the observations dY(t) by N Kalman filters. By taking advantage of the structure of the models in § 2, one can obtain an algorithm where the observations are processed by a *single Kalman filter*. To do so, we denote the model of the state process before incorporating predictive information as

$$dx(t) = A(t)x(t) dt + dU(t)$$
(35)

The hypothesis  $H_i$  corresponds to an added observation of the form

$$m_i = x(t_i) + v_i, \quad v_i \sim N(0, \pi_i)$$
 (36)

Hence the estimate  $\hat{x}_i(t)$  is given by

$$\hat{x}_i(t) = E\{x(t)|Y_t, m_i\}$$

Consequently,  $\hat{x}_i(t)$  can be viewed as a smoothed estimate that can be obtained by using the two-filter procedure described by Mayne (1966), Fraser (1967) and Wall et al. (1981) among others. The forward filter processes the information  $Y_t$  as

$$d\hat{x}(t) = A(t)\hat{x}(t) dt + \sum_{t} (t)C'(t)R^{-1}(t) [dy(t) - C(t)\hat{x}(t) dt]$$
(37)

$$\frac{d}{dt}\Sigma(t) = A(t)\Sigma(t) + \Sigma(t)A'(t) + Q(t)$$

$$- \Sigma(t)C'(t)R^{-1}(t)C(t)\Sigma(t)$$

$$\hat{x}(0) = m_0, \quad \Sigma(0) = \pi_0$$
(38)

where the matrices  $C_i(t)$  are assumed to be independent of i. Note that the forward filter is the same for all hypotheses. The backward filter for  $H_i$  is given by

$$\frac{d}{dt}\bar{x_i}(t) = A(t)\bar{x_i}(t) \tag{39}$$

$$\frac{d}{dt}\bar{\Sigma}_i(t) = A(t)\bar{\Sigma}_i(t) + \bar{\Sigma}_i(t)A'(t) - Q(t)$$
(40)

with the terminal conditions

$$\bar{x}_i(t_i) = m_i, \quad \bar{\Sigma}_i(t_i) = \pi_i$$
 (41)

so that we can identify  $\bar{x_i}(t) = \phi(t, t_i)m_i$  and  $\bar{\Sigma}_i(t) = \pi(t, t_i)$  which can be precomputed off-line, since they require no processing of measurements. The forward and backward estimates are combined to yield

$$\hat{x}_i(t) = \left[ \Sigma^{-1}(t) + \bar{\Sigma}_i^{-1}(t) \right]^{-1} \left[ \Sigma^{-1}(t)\hat{x}(t) + \bar{\Sigma}_i^{-1}(t)\bar{x}_i(t) \right]$$
(42)

where  $v_i$  is assumed independent of the Wiener processes U(t) and V(t).

The equation (42) can be substituted in (33) to obtain a recursive expression for  $P_i(t)$  driven by the observation dy(t) and the forward filter estimates  $\hat{x}(t)$ . Hence, the likelihood ratio given by eqn. (31) can be obtained by processing the observations with a single Kalman filter. This implementation has the advantage that the tracking filter produces the same estimate  $\hat{x}(t)$  whether or not there are predictive hypotheses, until a specific hypothesis is accepted, whereupon a new single model can be produced for the tracking algorithm.

#### 4. Detection of branching times

In § 2 and 3 we have studied the problem of selecting a tracking model for a trajectory, given some predictive information on the possible destination. However, these results were assuming a known origin, i.e. a known branching time. In practice, this branching time is seldom known, as ships are likely to change their mission or their destination while at sea. In this section, we develop similar hypothesis-testing procedures for processes where the branching times are not known. The basic approach is motivated by previous work by Willsky and Jones (1978) for detection of jumps in linear systems using a generalized maximum likelihood technique.

Assume that the original trajectory is modelled by

$$H_0: dx(t) = [A_0(t)x(t) + s_0(t)] dt + B(t) dU(t)$$
(43)

and that at an unknown time  $\theta$ , the ship changes destination to one of several locations described by Fig. 2, yielding N possible models described by eqns. (27), (28) and (29).

We assume that the observations dy(t) have been stored over a time interval [0, t], as well as the state estimate  $\hat{x}_0(t)$ . The problem consists of identifying two parameters: the destination i and the switch time  $\theta$ . We define the likelihood ratios for a given  $\theta$  as

$$L_{ij}(t; \theta) = \frac{dP(Y_t | \theta, H_i)}{dP(Y_t | \theta, H_i)}$$
(44)

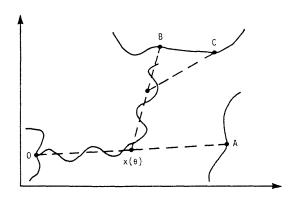


Figure 2. Detection of manoeuvres.

For a given  $\theta$ , (44) can be expressed as

$$L_{ij}(t;\theta) = \frac{P(H_i|\theta, Y_t)}{P(H_i|\theta)} \frac{P(H_j|\theta)}{P(H_i|\theta, Y_t)}$$

where the conditional probabilities  $P(H_i|\theta)$  are specified a priori and

$$P(H_i|\theta, Y_t) \triangleq P_i(t;\theta) \tag{45}$$

Separate  $Y_t$  into  $Y_{\theta}$  and  $Y[\theta, t]$  where

$$Y[\theta, t] = \sigma\{y(s), \theta \le s \le t\}$$

Then, one can process  $Y_{\theta}$  to get an initial distribution for  $x(\theta)$  as

$$P(x(\theta)|Y_{\theta}) = N(\hat{x}_{0}(\theta), \Sigma_{0}(\theta)) \tag{46}$$

The statistics (46) can now be used as initial conditions in (32)–(34) with the initial conditions

$$P_i(\theta; \theta) = P(H_i|\theta)$$

and the information  $Y[\theta, t]$  to obtain  $P_i(t; \theta)$ ,  $0 \le i \le N$ . The likelihood ratios can be obtained as in (31), so that

$$L_{i0}(t;\theta) = \exp\frac{1}{2} \int_{\theta}^{t} \left[ C(s)\hat{x}_{i}(s) - C(s)\hat{x}_{0}(s) \right] R^{-1}(s)$$

$$\left[ 2 \ dy(s) - C(s) \left[ \hat{x}_{i}(s) + \hat{x}_{0}(s) \right] \ ds \right]$$
(47)

which can be implemented by processing the information  $Y[\theta, t]$  using only one Kalman filter, as indicated in eqns. (37)–(42). Furthermore, this filter is identical for all values of  $\theta$ . This implies that the estimates  $\hat{x}_i(t)$  are independent of the branching time  $\theta$ .

However, our evolution equations for the identification probabilities  $P_i(r; \theta)$ , as well as the likelihood ratio  $L_{i0}(r; \theta)$ , depend on the value of the parameter  $\theta$ . This suggests a generalized maximum likelihood identification method, where

$$L_0^*(t; \theta) = \max_{0 \le i \le N} L_{i0}(t; \theta)$$
 (48)

and

$$L_0^*(t) = \max_{0 \le \theta \le t} L_0^*(t; \theta) \tag{49}$$

Other detection rules can be based on observation of the residuals. We have exhibited a simple detection procedure which can be implemented in an efficient manner. Note that the only reprocessing of information is carried out when evaluating the likelihood ratios. Once a destination has been identified, the new estimate of the state can be obtained from (42) without reprocessing the observations. Furthermore, if the optimal selection of models is independent of  $\theta$  over the range of interest, then the state estimate  $\hat{x}(t)$  can be obtained independently of eqn. (49). Thus, our implementation has the advantage that all the processing of information is carried out by a single Kalman filter, and while the identification procedure is taking place, this filter continues to produce state estimates which can be used for decision purposes.

### 5. A simple example

Consider the following two-dimensional Markov process:

$$x(t) = w_1(t)$$

$$y(t) = w_2(t)$$

where  $w_1$  and  $w_2$  are standard independent brownian motions. Consider the following predictive information: at time t = 14, the state x(t) will be distributed as

$$N\left(\begin{bmatrix} \pm 10 \\ \pm 10 \end{bmatrix}, \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}\right)$$

with equal probability. The situation is shown in Fig. 3.

(-10,-10)

(4)

Assume that, real-time measurements are taken, according to the equation

$$dz_1(t) = x_1(t) dt + dv_1(t)$$

$$dz_2(t) = y(t) dt + dv_2(t)$$

where  $v_1, v_2, w_1$  and  $w_2$  are all independent brownian motions. In this example, the equations (39) and (40) decouple in each coordinate, yielding

Figure 3. Possible destinations for Example 1.

(10,-10)

(3)

Furthermore, the Ricatti equation (38) decouples into two identical equations of the form

$$\frac{d}{dt}\,\Sigma(t) = -\,\Sigma^2 + 1$$

$$\Sigma(0) = 0$$

Table 1 shows the evolution of the probabilities when the true destination is (10, 10), and the trajectory from the origin to (10, 10), is a straight line of constant speed, using a sample of random numbers. The column  $\hat{x}$  indicates the estimated position using no predictive information. The column  $\hat{x}_s$  is the optimal mean, using the predictive information that all four estimations are equally likely. The actual random trajectory is indicated in the column x.

Time	$P_1$	$P_2$	$P_3$	$P_4$	x	у	â	ŷ	$\hat{X}_s$	$\hat{y}_s$
0.6	0.03	0.43	0.51	0.03	-0.37	0.43	-0.02	0.38	-0.04	0.92
1.2	0	0.75	0.25	0	0.93	0.15	0.14	0.76	0.48	1.38
1.8	0.05	0.95	0	0	1.0	0.14	0.74	0.37	1.4	1.0
2.4	0.08	0.92	0	0	2.7	0.23	1.5	0.29	2.2	0.96
3.0	0.01	0.99	0	0	2.5	1.8	2.5	0.56	3.0	1.26
3.6	0	1	0	0	2.3	2.6	2.5	0.97	3.1	1.7
5.0	0	1	0	0	2.4	4.93	1.9	4	2.6	4.6

Table 1.

Note the relatively quick identification of the destination, even though the local properties of brownian motion make the destination difficult to identify. In addition, note the improvement in performance of the filter with predictive information. Even though the average of the *a priori* predictive information corresponds to no information at all, the use of a recursive hypothesis estimation logic resulted in correct a posteriori identification which improved the estimator.

The results of Table 1 indicate how the predictive information is useful when the actual mean trajectory is a straight line to the destination. Assume now that the actual model for the trajectory is

$$x_{1}(t) = \frac{\sqrt{2}}{2} \left( t - \sin \frac{t\pi}{7} \right) + w_{1}(t)$$

$$x_{2}(t) = \frac{\sqrt{2}}{2} \left( \left( t + \sin \frac{t\pi}{7} \right) + w_{2}(t) \right)$$
(50)

This trajectory is depicted in Fig. 4.

Table 2 shows the resulting evolutions of the probabilities, the sample path of the process, and the filtered estimates with and without predictive information.

Although the identification of the correct destination occurs promptly, the results in Table 2 suggest that incorporating the correct predictive destination does not improve the performance of the estimator as much as it did in Table 1. This is the result of the mismatch between the linear models, incorporating predictive information, and the non-linear trajectories generated by (50).

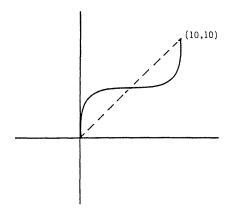


Figure 4. Mean ship trajectory for Table 2.

Time	$P_1$	$P_2$	$P_3$	$P_4$	х	у	â	ŷ	$\hat{X}_s$	$\hat{\mathcal{Y}}_s$
0.6	0.03	0.19	0.66	0.12	-0.38	-1.6	-0.17	0.24	-0.45	0.83
1.2	0.5	0.88	0.12	0	1.17	1.13	0.28	0.73	1.09	1.35
1.8	0	0.95	0.05	0	1.56	0.65	0.38	0.86	1.06	1.5
2.4	0	1	0	0	2.1	3.7	1.1	1.4	1.7	2.0
3.0	0	1	0	0	0.5	3.3	1.3	2.1	2.0	2.78
5.0	0	1	0	0	3.3	5.4	2.1	4.1	2.8	4.6
7.0	0	1	0	0	7.0	4.2	4.7	4.6	5.3	5.2
8.0	0	1	0	0	5.3	5.8	5.5	4-4	6-1	5.1

Table 2.

# 6. Conclusions

This paper has shown how some information about the possible destinations of a ship can be incorporated in the modelling of its trajectories. These results are potentially useful for the tracking of multiple objects (cf. Bar-Shalom (1978) and Morefield (1977)). One possible approach is to generate multiple models based on the likely associations of data and use a sequential evaluation of these models as described in § 3. In addition, when a region is congested with traffic, the additional knowledge of the destination of the various ships in the area is very valuable for the track reconstruction problem (one need only match the outgoing tracks with the ships' destinations).

Another appealing feature of the algorithms presented here is their capacity to evaluate the validity of predictive information. In particular, §§ 3 and 4 discuss algorithms where the predictive information is not used in the tracking filter until a likelihood-ratio test confirms its validity.

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