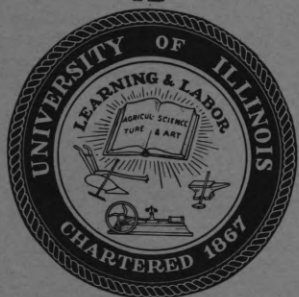




Coordinated
Science
Laboratory



UNIVERSITY OF ILLINOIS - URBANA, ILLINOIS

ON MINIMAL-VARIANCE CONTROL
OF LINEAR SYSTEMS
WITH QUADRATIC LOSS

Michael K. Sain

R-240

JANUARY 1965

COORDINATED SCIENCE LABORATORY
UNIVERSITY OF ILLINOIS
URBANA, ILLINOIS

Contract DA-28-043-AMC-00073(E)

The research reported in this document was made possible by support extended to the University of Illinois, Coordinated Science Laboratory, jointly by the Department of the Army, Department of the Navy (Office of Naval Research), and the Department of the Air Force (Office of Scientific Research) under Department of Army Contract DA-28-043-AMC-00073(E).

ERRATA

Page 22, Eqn. (3.47) $\int_{t_0}^T$ instead of \int_0^T .

Page 23, Eqn. (3.52) same correction as Eqn. (3.47).

Page 32, Eqn. (5.2) $\Phi^T(\tau-\sigma)$ instead of $\Phi^T(t-\sigma)$.

ON MINIMAL-VARIANCE CONTROL OF LINEAR
SYSTEMS WITH QUADRATIC LOSS

Michael Kent Sain, Ph.D.
Department of Electrical Engineering
University of Illinois, 1965

ABSTRACT

The problem of minimizing the variance of a quadratic performance index, in the presence of control noise whose properties are known a priori, has been studied for linear, constant systems with open-loop control. The expected value of the performance index is constrained to be a positive number, which cannot be less than the optimal mean index when variance is free.

Solution is by means of the calculus of variations, which is applied to an equivalent noise-free problem. The necessary (Euler) equations are integro-differential and have a kernel matrix derived from output correlation functions. In general, these equations contain a forcing vector which depends upon the third-moment properties of the noise process. For systems in which the state vector can be chosen as the output, the existence of an inverse for the kernel matrix can be related to the total state controllability of an equivalent linear, noise-free plant which incorporates statistical data from the disturbance process.

A maximum estimate for the number of eigenvalues is given, and is refined for the special case of single-input, state-output control. Necessary and sufficient conditions for a unique solution can be found in specific examples.

ACKNOWLEDGMENT

The author sincerely appreciates the guidance and suggestions of Professor J. B. Cruz, Jr. and many helpful discussions with Professors M. E. Van Valkenburg, W. R. Perkins, and R. A. Rohrer. He wishes also to thank all those who, as teachers and friends, have made these results possible.

In addition to the assistantship support of the Coordinated Science Laboratory of the University of Illinois, the author acknowledges the fellowship support of the National Electronics Conference and the National Science Foundation.
National Science Foundation.

TABLE OF CONTENTS

	Page
LIST OF ILLUSTRATIONS	v
1. INTRODUCTION	1
2. PROBLEM STATEMENT	4
3. OPTIMIZATION.	7
3.1 A Noise-Free Equivalent Problem	7
3.2 Necessary Conditions and Equations	9
3.3 First-Order Examples	12
3.3.1 A Constant Plant	12
3.3.2 A Time-Varying Plant	22
3.4 Relation to Minimal-Expectation Problem	23
4. SINGLE-INPUT CONTROL.	25
4.1 Structure of Solution	25
4.2 A Degenerate Example	30
5. MULTI-INPUT SYSTEMS	32
5.1 A Typical Form for $Q(t,\tau)$	32
5.2 A Method of Analysis	34
5.3 Second-Order Example	37
6. CONCLUSIONS	41
6.1 Summary	41
6.2 Problems for Further Study	42
BIBLIOGRAPHY	43
VITA	45

LIST OF ILLUSTRATIONS

Figure	Page
1. Variance J° and Mean K_{\circ} as a Function of $\mu(\mu > 0)$	16
2. Variance J° and Mean K_{\circ} as a Function of $\mu(\mu < 0)$	17
3. Mean K_{\circ} as a Function of Variance J°	19
4. Variance J° and Mean K_{\circ} as a Function of $\mu(\mu > 0)$ with Parameter $z(0)$	20
5. Variance J° and Mean K_{\circ} as a Function of Duration T with Parameter a, for Constant μ	21

1. INTRODUCTION

A typical problem in the theory of optimal control may be defined in the following manner. Select an m -dimensional control vector $\underline{u}(t)$, $t \in [0, T]$, such that the linear, constant system

$$\dot{\underline{x}} = A\underline{x} + B\underline{u} , \quad (1.1)$$

$$\underline{y} = H\underline{x} \quad (1.2)$$

moves from an initial output \underline{y}_0 to a final output \underline{y}_T while minimizing the performance index

$$J = \int_0^T (\underline{y}^T P \underline{y} + \underline{u}^T R \underline{u}) dt. \quad (1.3)$$

The vectors \underline{x} and \underline{y} have dimension n and r ; the time-invariant matrices A , B , H , P , and R are $n \times n$, $n \times m$, $r \times n$, $r \times r$, and $m \times m$ respectively; P and R are symmetric and positive-definite.

It often happens that the control input is accompanied by undesirable disturbances. For example, an antenna positioning device may experience fluctuating loads because of wind gusts. One method of representing such effects is to add a term $C\underline{y}$ in the right member of (1.1), where C is an $n \times p$ constant matrix and \underline{y} is a p -vector, continuous parameter random process.

In such cases, it is possible to replace (1.3) with an essentially different criterion, such as the probability of achieving a given task within prescribed tolerances. An example related to this viewpoint has been formulated by Pontryagin, et. al. [1]. However, if the random processes associated with an optimization problem are true disturbances --

that is, not completely dominating system behavior, yet adding uncertainty to any performance -- then the noise-free index (1.3) appears to retain its original significance. It is itself a random variable, assuming values corresponding to sample functions of the \underline{y} process.

From this latter point of view, the statistical moments of J may be used to define extensions of (1.3). If the system is to be used for a large number of identical performance attempts, as for example at a target practice, a criterion of the form $E\{J\}$ tends to give satisfactory results. Thus it is reasonable to believe that a control which optimizes the mean or ensemble average $E\{J\}$ will, in at least one instance, perform very nearly as expected. This minimal-expectation problem has received considerable attention in the literature. Kushner [2] has developed a stochastic maximum principle for these indices; Kipiniak [3] has investigated the class with the calculus of variations; Bellman and others [4,5,6,7] have used dynamic programming.

Consider, however, an application in which but a single performance is possible; only one opportunity is available to achieve the design objectives. Average behavior may then be quite unsatisfactory, especially if the standard deviation is large. With this motivation, a control will be optimal in the sense of this investigation if it brings about a local minimum of index variance, while satisfying a constraint on index expectation. Very little explicit work on this topic has been published. For systems of the type (1.1), with a term $C\underline{y}$ included, problems which seek to minimize the mean value of a quadratic form in final state or output seem close in concept to the work of this paper. But it can be shown, as

Dreyfus [8] has demonstrated, that these reduce to the solution of the same system with no disturbances. Moreover, as observed by Kushner [9], problems which attempt to optimize $E\{J\}$, J being defined as in (1.3), have this property in general. This does not occur in the minimal-variance problem.

As considered herein, the minimal-variance control problem can, of course, be written in final-value form, as in Orford [10]; but the system equations are then nonlinear, so that many analytical advantages are lost. Equations in this paper are integro-differential, but linear nonetheless. Finally, it is also possible to make use of the terminology of minimal-variance, unbiased-estimation theory. Because the functional estimator is quadratic in the unknowns, however, the problem differs from the work of other investigators, for example, Kalman [11].

2. PROBLEM STATEMENT

In this initial investigation, no attempt has been made to formulate the minimal-variance problem in its most general form. A high priority on analytical results which display the structure of solution makes certain simplifying assumptions necessary.

This analysis therefore considers a linear system, additive control noise, a quadratic performance index, a priori knowledge about the statistical nature of the disturbance process, and open-loop control. If the state is known as time evolves, a performance improvement accrues from application of the open-loop, optimal feedback realization [8].

The problem is to select a control vector $\underline{u}(t)$, $t \in [0, T]$, so that the system

$$\dot{\underline{x}} = \underline{A}\underline{x} + \underline{B}\underline{u} + \underline{C}\underline{v} , \quad (2.1)$$

$$\underline{y} = \underline{H}\underline{x} \quad (2.2)$$

moves from an expected initial output

$$E\{\underline{y}(0)\} = \underline{y}_0 \text{ (given)} \quad (2.3)$$

to an expected final output

$$E\{\underline{y}(T)\} = \underline{y}_T \text{ (given)} \quad (2.4)$$

at time T (given), while satisfying the constraint

$$E\{K\} = K_0 \text{ (given)}, \quad (2.5)$$

and achieving a local minimum J^0 of the functional

$$J = E\{(K - K_0)^2\} , \quad (2.6)$$

where

$$K \triangleq \int_0^T (\underline{y}^T P \underline{y} + \underline{u}^T R \underline{u}) dt. \quad (2.7)$$

In expressions (2.1) through (2.7), \underline{x} is an n -vector; \underline{u} is an m -vector; \underline{y} is an r -vector; and \underline{v} is a p -vector random process which characterizes the control disturbance. It follows that \underline{x} and \underline{y} are also random processes. The constant matrices A , B , C , and H have dimension $n \times n$, $n \times m$, $n \times p$, and $r \times n$ respectively; the positive-definite, symmetric matrices P and R have dimension $r \times r$ and $m \times m$. The initial plant state $\underline{x}(0)$ is a random vector independent of the \underline{v} process. In (2.5), K_0 is positive and bounded below by the solution of the minimal-expectation problem. The duration T of operation is finite.

The solution replaces the given problem by an equivalent deterministic counterpart, applies the calculus of variations, and obtains necessary (Euler) equations. Since (2.1) and (2.2) are linear, the method is feasible if \underline{v} has at least a stationary zero mean, and if it is permissible to exchange the order of expectation and integration in (2.5) and (2.6).

Although the necessary equations depend only upon the above assumptions, discussion of the solution simplifies if \underline{v} is also gaussian, white, and has mutually uncorrelated components. The assumption that \underline{v} is gaussian corresponds to many physical cases in which the noise process comes about as a result of the combined effects of a large number of minute disturbances. Moreover, Sivan [12] has shown, for linear plants with quadratic indices, that the gaussian case is "worst to control," in the minimal-expectation problem; that is, the index value is larger for gaussian than for other noises. This result indicates an additional incentive for the gaussian assumption. If \underline{v} is not white, it can be

represented by the output of a linear system driven by white noise [13].

The mathematical meaning of (2.1) can be related to the concept of a generalized random process [9,11]; or it may be discussed in terms of a differential stochastic process, as for example in Doob [14]. However, the solution can also proceed in the manner of this paper, since only integrals of the \underline{y} process occur. Such an approach, though formal insofar as the definition of (2.1) is concerned, may be put on a rigorous basis by means of its integral equivalent [11,15].

The necessary equations for time-varying systems are straightforward extensions of those for the constant system of this section. They appear in Section 3.3.2.

3. OPTIMIZATION

Because it is assumed that a priori information concerning the \underline{y} process is available, the averaging operations indicated in the previous section can be carried out directly to convert the problem into a disturbance-free counterpart. The calculus of variations is then applicable directly; and linear, integro-differential equations arise as necessary (Euler) conditions for solution. Besides the output controllability conditions implied in (2.3) and (2.4), similar requirements can be stated in controllability terms to establish the nonsingularity of the Euler kernel matrix, if H has an inverse. A constant and a time-varying example reflect the elementary characteristics of the minimal-variance control, and lead naturally to a discussion of relationships with the minimal-expectation problem.

3.1 A Noise-Free Equivalent Problem

Let \underline{y} be the sum of two parts \underline{y}_u and \underline{y}_v , which depend upon \underline{u} and \underline{v} respectively. Such a division is possible since (2.1) and (2.2) are linear. Thus

$$\underline{x} = \underline{z} + \underline{w} , \quad (3.1)$$

where
$$\dot{\underline{z}} = \underline{A}\underline{z} + \underline{B}\underline{u}, \quad \underline{z}(0) = E\{\underline{x}(0)\} , \quad (3.2)$$

and
$$\underline{w} = \int_0^t \Phi(t-\sigma) \underline{C}\underline{v}(\sigma) d\sigma . \quad (3.3)$$

In (3.3), $\Phi(t)$ is the state transition matrix [16] corresponding to (2.1).

It follows immediately that

$$\underline{y}_u = \underline{H}\underline{z}; \quad \underline{y}_v = \underline{H}\underline{w} ; \quad (3.4)$$

and
$$E\{\underline{y}_v\} = \underline{0} . \quad (3.5)$$

Interchange of expectation and integration in (2.5) gives

$$K_o = \int_0^T (E\{\underline{y}^T P \underline{y}\} + \underline{u}^T R \underline{u}) dt , \quad (3.6)$$

since \underline{u} is an open-loop control. Moreover, (2.6) and (2.7) lead in a similar manner to

$$J = \int_0^T \int_0^T E\{\Delta_{yp}(t) \Delta_{yp}(\tau)\} dt d\tau , \quad (3.7)$$

where

$$\Delta_{yp} = \underline{y}^T P \underline{y} - E\{\underline{y}^T P \underline{y}\} . \quad (3.8)$$

Now $\underline{y} = \underline{y}_u + \underline{y}_v$, and

$$E\{\underline{y}^T P \underline{y}\} = \underline{y}_u^T P \underline{y}_u + E\{\underline{y}_v^T P \underline{y}_v\} . \quad (3.9)$$

An equivalent problem, therefore, is to select a control vector $\underline{u}(t)$, $t \in [0, T]$, so that the system

$$\dot{\underline{z}} = A \underline{z} + B \underline{u} \quad (3.10)$$

satisfies the boundary conditions

$$H \underline{z}(0) = \underline{y}_0, \quad H \underline{z}(T) = \underline{y}_T , \quad (3.11)$$

obeys the constraint

$$\int_0^T (\underline{z}^T H^T P H \underline{z} + \underline{u}^T R \underline{u} + \sqrt{q_1(t, t)}) dt = K_o , \quad (3.12)$$

and minimizes the index

$$J = \int_0^T \int_0^T [4 \underline{z}^T(\tau) H^T Q(\tau, t) H \underline{z}(t) + 4 \underline{z}^T(t) H^T Q(t, \tau) + q_2(\tau, t) - q_1(\tau, t)] d\tau dt . \quad (3.13)$$

It has been convenient to define

$$Q(\tau, t) = PE\{y_v(\tau)y_v^T(t)\}P, \quad (3.14)$$

$$g(\tau, t) = PE\{y_v(\tau)[y_v^T(t)Py_v(t)]\}, \quad (3.15)$$

$$q_1(\tau, t) = E\{y_v^T(\tau)Py_v(\tau)\}E\{y_v^T(t)Py_v(t)\}, \quad (3.16)$$

$$q_2(\tau, t) = E\{[y_v^T(\tau)Py_v(\tau)][y_v^T(t)Py_v(t)]\}. \quad (3.17)$$

It will be seen in the following section that the scalar functions $q_1(t, \tau)$ and $q_2(t, \tau)$ affect K_0 and J^0 but do not appear in the necessary equations. The matrix $Q(t, \tau)$ and the vector $g(t, \tau)$ contribute to the kernel and driving functions respectively in these relations.

3.2 Necessary Conditions and Equations

In the calculus of variations, (3.12) is an isoperimetric constraint.

The technique of solution appends (3.10) and (3.12) to (3.13) with Lagrange multipliers $\underline{\lambda}$ and μ to give a new, unconstrained index

$$\begin{aligned} J' = & \int_0^T [\underline{\lambda}^T(\dot{\underline{z}} - A\underline{z} - B\underline{u}) + \mu(\underline{z}^T H^T P H \underline{z} + \underline{u}^T R \underline{u} + \sqrt{q_1(t, t)}) \\ & + \int_0^T (4\underline{z}^T(\tau) H^T Q(\tau, t) H \underline{z}(t) + 4\underline{z}^T(t) H^T g(t, \tau) \\ & + q_2(\tau, t) - q_1(\tau, t)) d\tau] dt. \end{aligned} \quad (3.18)$$

The first variation is

$$\begin{aligned} \delta J' = & \underline{\lambda}^T(T) \delta \underline{z}(T) - \underline{\lambda}^T(0) \delta \underline{z}(0) + \int_0^T (2\mu \underline{u}^T R - \underline{\lambda}^T B) \delta \underline{u} dt \\ & + \int_0^T [2\mu \underline{z}^T H^T P H - \dot{\underline{\lambda}}^T - \underline{\lambda}^T A \\ & + 4 \int_0^T (2\underline{z}^T(\tau) H^T Q(\tau, t) H + g^T(t, \tau) H) d\tau] \delta \underline{z}(t) dt. \end{aligned} \quad (3.19)$$

The necessary equations are then

$$2\mu R\underline{u} = B^T \underline{\lambda} ; \quad (3.20)$$

$$4H^T \int_0^T (2Q(t,\tau)H\underline{z}(\tau) + \underline{q}(t,\tau))d\tau + 2\mu H^T P H \underline{z} = \dot{\underline{\lambda}} + A^T \underline{\lambda} . \quad (3.21)$$

The foregoing procedure uses the property that

$$Q^T(t,\tau) = Q(\tau,t) , \quad (3.22)$$

which follows from (3.14).

Consider now the special case in which \underline{v} is white and wide-sense stationary [14]; then

$$E\{\underline{v}(t)\underline{v}^T(\tau)\} = S \delta(t-\tau) ; \quad (3.23)$$

the covariance matrix S is symmetric and positive semi-definite. Thus there exists a matrix M with the property that

$$S = M M^T . \quad (3.24)$$

The kernel matrix $H^T Q(t,\tau)H$ of (3.21) can be written, by means of (3.3), (3.14), and (3.24), as

$$H^T Q(t,\tau)H = H^T P H \int_0^{(t,\tau)} \Phi(t-\sigma) C M M^T C^T \Phi^T(\tau-\sigma) d\sigma H^T P H , \quad (3.25)$$

where the notation (\cdot, \cdot) selects the smaller of its arguments. If H^{-1} exists (in which event the output relations can be incorporated into the plant equations), the kernel matrix is nonsingular if and only if the Gramian matrix

$$P((t,\tau)) = \int_0^{(t,\tau)} \Phi((t,\tau)-\sigma) C M M^T C^T \Phi^T((t,\tau)-\sigma) d\sigma \quad (3.26)$$

has an inverse for all $(t,\tau) \in [0,T]$. This requirement is equivalent to

the condition [17] that the plant (3.10), with B replaced by CM, be totally state controllable on the interval, which fact can be determined by the test

$$\text{rank} \begin{bmatrix} \text{CM} & \vdots & \text{ACM} & \vdots & \text{A}^2\text{CM} & \vdots & \dots & \vdots & \text{A}^{n-1}\text{CM} \end{bmatrix} = n. \quad (3.27)$$

This interpretation is not possible if H does not have an inverse.

Another necessary condition, valid for any H, is that

$$\text{rank} \begin{bmatrix} \text{HB} & \vdots & \text{HAB} & \vdots & \text{HA}^2\text{B} & \vdots & \dots & \vdots & \text{HA}^{n-1}\text{B} \end{bmatrix} = r, \quad (3.28)$$

which implies total output controllability of (3.10) with output equations $\underline{y}_u = \text{H}\underline{z}$ and assures that (3.11) can be satisfied for at least one control \underline{u} .

The vector $\underline{q}(t, \tau)$ derives from the third-moment properties of the \underline{y} process; and so the optimal control must use, in general, more information than that provided by the covariance matrix of the disturbance. However, if \underline{y} is gaussian, white, and has mutually uncorrelated components, then $\underline{q}(t, \tau)$ vanishes; and the equations (3.20) and (3.21) are homogeneous.

Consider also the second variation

$$\begin{aligned} \delta^2 J' = & 2\mu \int_0^T \delta \underline{z}^T \text{H}^T \text{P} \text{H} \delta \underline{z} dt + 2\mu \int_0^T \delta \underline{u}^T \text{R} \delta \underline{u} dt \\ & + 8 \int_0^T \int_0^T \delta \underline{z}^T(\tau) \text{H}^T \text{Q}(\tau, t) \text{H} \delta \underline{z}(t) d\tau dt. \end{aligned} \quad (3.29)$$

The matrices P and R are positive-definite; and the third term is non-negative; thus $\mu > 0$ is sufficient to ensure (3.29) positive, and thereby that the solution is a local minimum. Although this technique may not correspond to classical sufficiency conditions [18], it has fundamental implications in the following example.

3.3 First-Order Examples

In order to demonstrate the solution of (3.10), (3.20), and (3.21), it is instructive to examine the first-order ($n=1$) case. A detailed investigation follows for a constant system; and some corresponding properties are shown for a time-varying example.

3.3.1 A Constant Plant

Let the system equations be

$$\dot{x} = -ax + u + v, \quad (3.30)$$

$$y = x, \quad (3.31)$$

where $a > 0$ and

$$\varphi(t) = e^{-at}. \quad (3.32)$$

If $P = R = 1$, then

$$Q(t, \tau) = \frac{\sigma_v^2}{2a} (e^{-a|t-\tau|} - e^{-a(t+\tau)}); \quad (3.33)$$

$$q(t, \tau) = \frac{\alpha_v^3 e^{-a(t+2\tau)}}{3a} (e^{3a(t, \tau)} - 1); \quad (3.34)$$

$$q_1(t, \tau) = \frac{(\sigma_v^2)^2}{4a^2} (e^{-2at} - 1)(e^{-2a\tau} - 1); \quad (3.35)$$

$$\begin{aligned} q_2(t, \tau) &= \frac{\beta_v^4}{4a} (e^{-2a|t-\tau|} - e^{-2a(t+\tau)}) \\ &+ \frac{(\sigma_v^2)^2}{4a^2} (e^{-2at} - 1)(e^{-2a\tau} - 1) \\ &+ \frac{(\sigma_v^2)^2}{2a^2} (e^{-2a(t+\tau)} - 2e^{-2a[t, \tau]_+} e^{-2a|t-\tau|}), \end{aligned} \quad (3.36)$$

where $[\cdot, \cdot]$ selects the larger of its arguments. Here v is white and

stationary, with second, third, and fourth moments σ_v^2 , α_v^3 , and β_v^4 respectively.

Because there is only one control variable, (3.10), (3.20), and (3.21) combine to give a single integro-differential equation,

$$\begin{aligned} & \frac{2\sigma_v^2}{a} \int_0^T (e^{-a|t-\tau|} - e^{-a(t+\tau)}) z(\tau) d\tau + \frac{\alpha_v^3}{a^2} (1 - e^{-at}) \\ & + \frac{\alpha_v^3}{3a^2} e^{-2aT} (e^{-at} - e^{2at}) + \mu z = \mu z^{(2)} - a^2 z. \end{aligned} \quad (3.37)$$

The superscript notation (i) indicates the ith time derivative. Equation (3.37) is equivalent to an ordinary differential equation of fourth order. The reader may verify that

$$\begin{aligned} z^{(4)} - (1+2a^2)z^{(2)} + \left[\frac{4\sigma_v^2}{\mu} + a^2(a^2+1) \right] z \\ = - \frac{\alpha_v^3}{\mu} [1 + e^{-2a(T-t)}] \end{aligned} \quad (3.38)$$

is equivalent to (3.37) by a twofold time differentiation and substitution.

The solution is

$$\begin{aligned} z(t) = \sum_{i=1}^4 c_i e^{p_i t} - \frac{\alpha_v^3}{4\sigma_v^2 + \mu a^2(a^2+1)} \\ - \frac{\alpha_v^3 e^{-2a(T-t)}}{4\sigma_v^2 + 3\mu a^2(3a^2-1)}, \end{aligned} \quad (3.39)$$

where the p_i are the (distinct) roots of

$$p_i^4 - (1+2a^2)p_i^2 + \frac{4\sigma_v^2}{\mu} + a^2(a^2+1) = 0. \quad (3.40)$$

In case $\mu = 16\sigma_v^2$ or $\mu = -4\sigma_v^2/[a^2(a^2+1)]$, the p_i are not distinct; and appropriate changes occur in (3.39).

Two conditions arise from (3.11); and these are

$$\sum_{i=1}^4 c_i = z(0) + \frac{\alpha_v^3}{4\sigma_v^2 + \mu a^2(a^2+1)} + \frac{\alpha_v^3 e^{-2aT}}{4\sigma_v^2 + 3\mu a^2(3a^2-1)} ; \quad (3.41)$$

$$\sum_{i=1}^4 c_i e^{p_i T} = z(T) + \frac{\alpha_v^3}{4\sigma_v^2 + \mu a^2(a^2+1)} + \frac{\alpha_v^3}{4\sigma_v^2 + 3\mu a^2(3a^2-1)} \quad (3.42)$$

Two more conditions result from substituting (3.39) into (3.37).

Linearly independent functions of the form $e^{p_i t}$ and $e^{\pm at}$ result from this substitution. Since the coefficients of $e^{p_i t}$ terms vanish identically if p_i is a solution of (3.40), a necessary and sufficient condition for (3.39) to satisfy (3.37) is that the coefficients of $e^{\pm at}$ terms also vanish identically. This will occur if the following two conditions are satisfied:

$$\begin{aligned}
& \sum_{i=1}^4 c_i \left[\frac{e^{(p_i - a)T}}{p_i - a} + \frac{2a}{a^2 - p_i^2} \right] - \frac{\alpha_v^3}{4\sigma_v^2 + \mu a^2 (a^2 + 1)} \left[\frac{2}{a} - \frac{e^{-aT}}{a} \right] \\
& - \frac{\alpha_v^3}{4\sigma_v^2 + 3\mu a^2 (3a^2 - 1)} \left[\frac{e^{-aT}}{a} - \frac{2e^{-2aT}}{3a} \right] \\
& + \frac{\alpha_v^3}{6a\sigma_v^2} (3 - e^{-2aT}) = 0 ;
\end{aligned} \tag{3.43}$$

$$\begin{aligned}
& \sum_{i=1}^4 c_i \left[\frac{e^{(p_i - a)T}}{p_i - a} \right] + \frac{\alpha_v^3}{4\sigma_v^2 + \mu a^2 (a^2 + 1)} \left[\frac{e^{-aT}}{a} \right] \\
& - \frac{\alpha_v^3}{4\sigma_v^2 + 3\mu a^2 (3a^2 - 1)} \left[\frac{e^{-aT}}{a} \right] = 0 .
\end{aligned} \tag{3.44}$$

These equations are solvable for the constants $\{c_i\}$, $i = 1, 2, 3, 4$; the control follows from (3.10).

If v is a gaussian process, then $\alpha_v^3 = 0$. Let $\sigma_v^2 = 1$, $z(0) = 1$, $z(T) = 0$, $a = 0.1$, and $T = 1$. Figures 1 and 2 present curves of K_0 and J^0 versus μ . For a given K_0 , it may happen that more than one value of μ can satisfy the constraint (3.12); however, the numerical values of J^0 are smallest for positive μ . It is evident that the sensitivity of J^0 and K_0 to changes in μ is uniformly better for $\mu > 0$.

In general, as $\mu \rightarrow 0^+$, the solution corresponds to a minimal-variance problem with no constraints on $E\{K\}$. The control is impulsive, since \underline{u} is not weighted in the J index (3.13). As $\mu \rightarrow \infty$, the solution

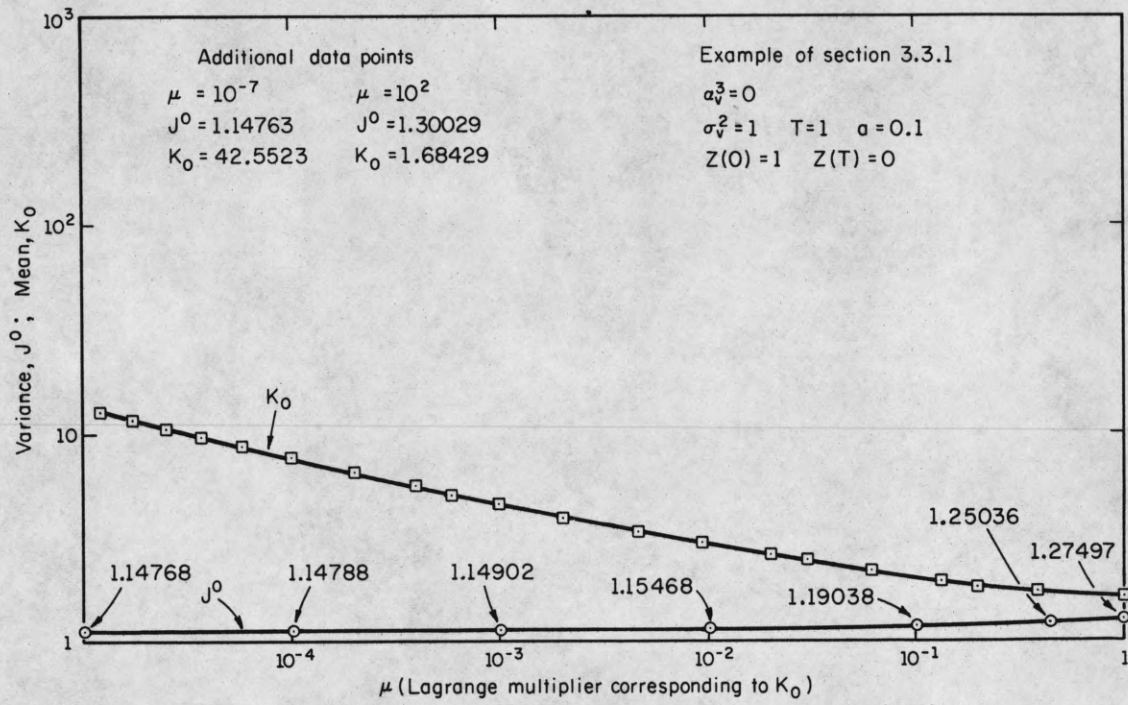


Fig. 1. Variance J^0 and Mean K_0 as a Function of μ ($\mu > 0$).

Fig. 2. Variance J° and Mean K_{\circ} as a Function of μ ($\mu < 0$).

minimizes $E\{K\}$, with no constraint on J ; Section 3.4 contains a discussion of this case.

By analysis of the c_i equations, (3.41) through (3.44), it is found that the excursions of Figure 2 are such as to satisfy a transcendental equation of the type

$$(a^2 - \alpha_1^2)(a + \beta_3 \cot \beta_3 T) = (a^2 + \beta_3^2)(a + \alpha_1 \coth \alpha_1 T), \quad (3.45)$$

where $p_i = \alpha_i + j\beta_i$. The character of (3.45) indicates an increasing number of excursions as $\mu \rightarrow 0^-$.

For positive μ , K_0 as a function of J^0 appears in Figure 3. The optimal variance J^0 changes relatively little as K_0 varies, since the contribution to J^0 by the process v dominates that of the control u (through z). Increased J^0 changes occur if $z(0)$ becomes larger in comparison with σ_v^2 . Of course, K_0 also grows in this case. Figure 4 substantiates these observations. The singular properties of u , as $\mu \rightarrow 0^+$, are also observed in these curves.

The relative size of J^0 and K_0 , for a fixed μ , depends upon T and a , if $z(0)$ and σ_v^2 are constant, as seen in Figure 5. Asymptotic behavior of K_0 and J^0 occurs because of steady v process contributions after controlled dynamic responses have decreased, which is in agreement with (3.40), an equation that does not contain T .

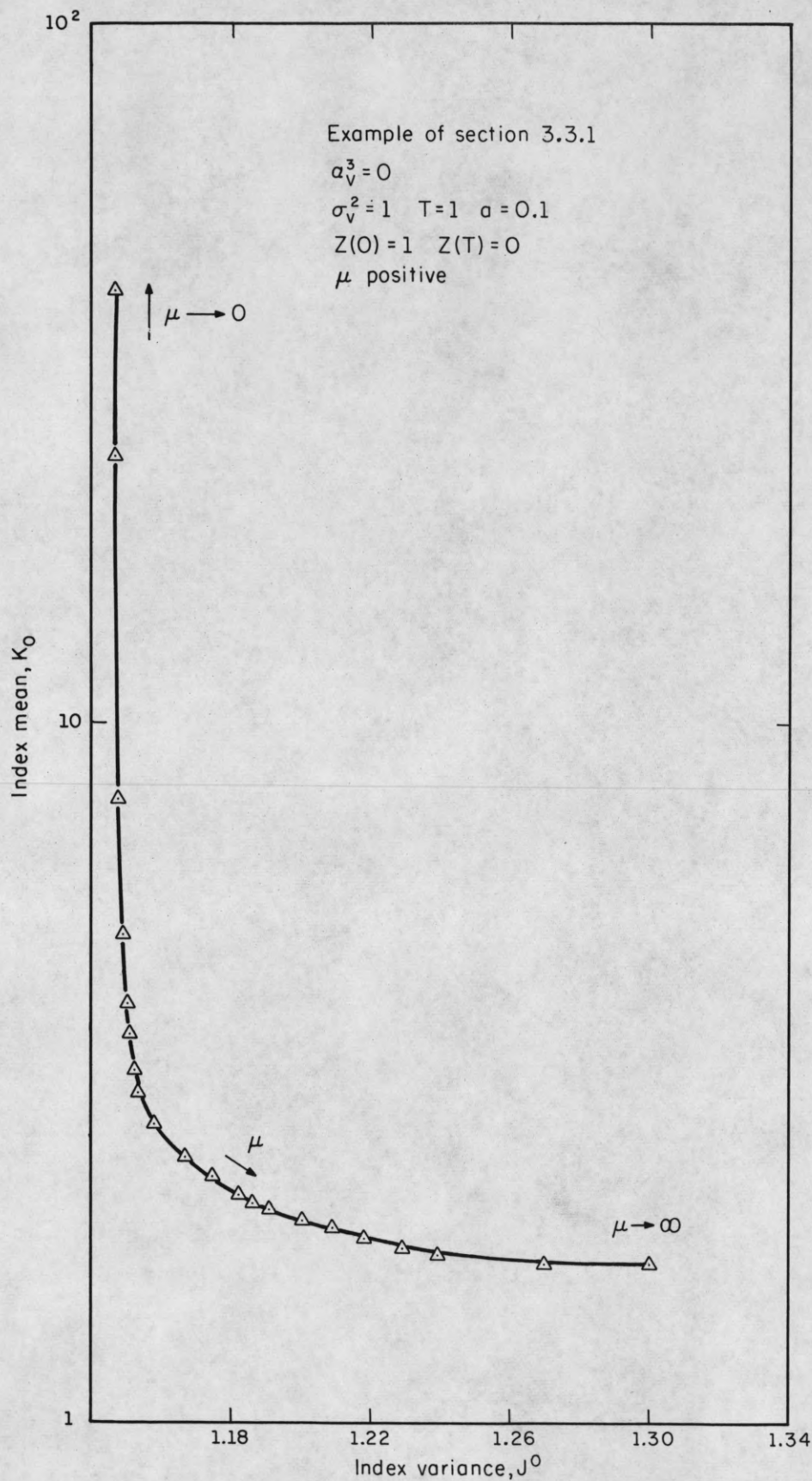


Fig. 3. Mean K_0 as a Function of Variance J^0 .

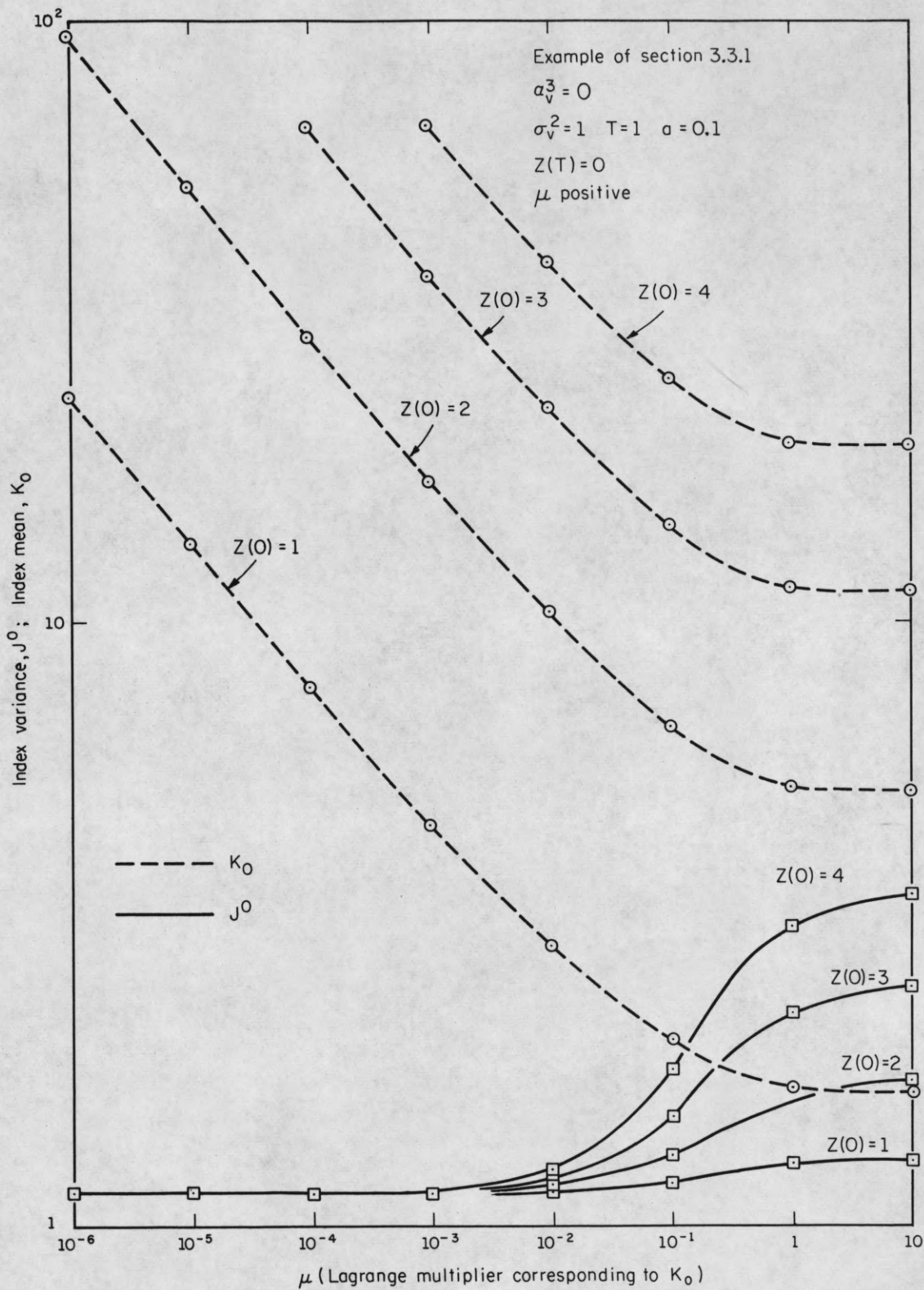


Fig. 4. Variance J^0 and Mean K_0 as a Function of μ ($\mu > 0$) with Parameter $z(0)$.

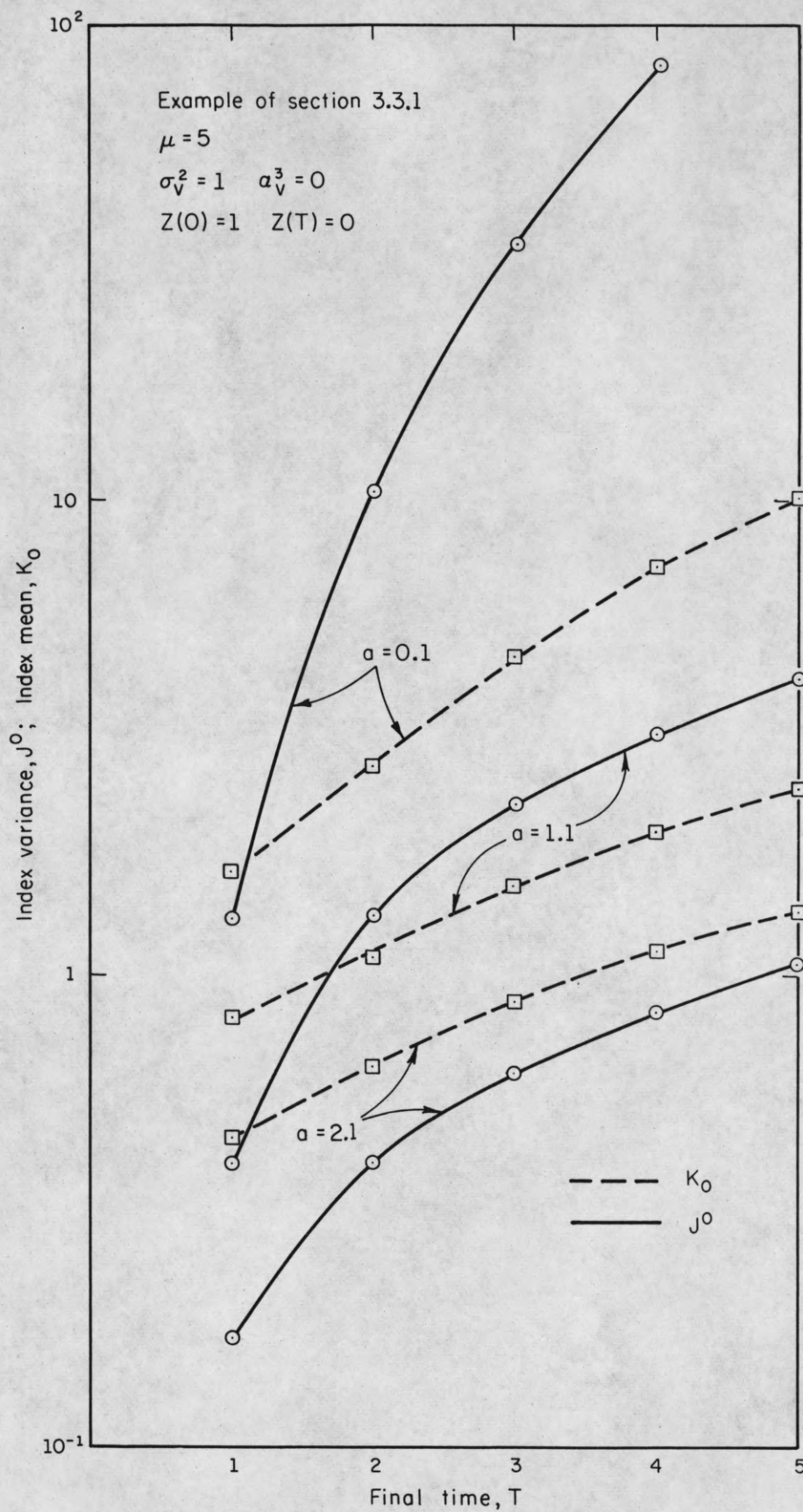


Fig. 5. Variance J^0 and Mean K_0 as a Function of Duration T with Parameter a , for Constant μ .

3.3.2 A Time-Varying Plant

In case the system is time-varying, the optimization proceeds again as in Section 3.2 to give the necessary equations

$$2\mu R(t)\underline{u} = B^T(t)\underline{\lambda} , \quad (3.46)$$

$$4H^T(t) \int_0^T [2Q(t,\tau)H(\tau)\underline{z}(\tau) + \underline{q}(t,\tau)]d\tau \\ + 2\mu H^T(t)P(t)H(t)\underline{z} = \dot{\underline{\lambda}} + A^T(t)\underline{\lambda} . \quad (3.47)$$

Corresponding modifications must be made in the plant equations, boundary conditions, and expressions for K and J; definitions (3.14) through (3.17) now depend on t_0 through (3.3), which must have a lower limit t_0 and incorporate a transition matrix $\Phi(t,\sigma)$.

Although many of the properties of the preceding example carry over to the time-varying example below, lack of a closed-form solution prevents development of a complete analogy. Let the system be

$$\dot{\underline{x}} = -\frac{1}{t} \underline{x} + \underline{u} + \underline{v} , \quad (3.48)$$

$$\underline{y} = \underline{x} , \quad (3.49)$$

where

$$\varphi(t,t_0) = t_0/t , \quad (3.50)$$

and

$$Q(t,\tau) = \frac{\sigma_v^2 P^2}{3t\tau} ((t,\tau)^3 - t_0^3) . \quad (3.51)$$

Here $P = p$, $R = r$, and v is white and gaussian. The integro-differential equation corresponding to (3.37) is

$$\frac{4\sigma_v^2 p^2}{3t} \int_0^T \frac{((t,\tau)^3 - t_0^3)}{\tau} z(\tau) d\tau$$

$$+ \mu p z = \mu r (z^{(2)} - \frac{2}{t^2} z) . \quad (3.52)$$

The method of differentiation and substitution is applicable again; and the ordinary differential equation corresponding to (3.38) is

$$z^{(4)} - \left[\frac{4}{t^2} + \frac{p}{r} \right] z^{(2)} + \frac{8}{t^3} z^{(1)}$$

$$+ \left[\frac{2}{r} \left(\frac{p}{t^2} + \frac{2\sigma_v^2 p^2}{\mu} \right) - \frac{8}{t^4} \right] z = 0 . \quad (3.53)$$

Thus the optimal trajectory consists once more of $4n$ characteristic functions. This property appears general, if \underline{u} has only one component. Section 4 considers this class of problems in more detail for constant systems.

3.4 Relation to Minimal-Expectation Problem

Consider the problem of selecting a control vector $\underline{u}(t)$, $t \in [0, T]$, so that the system (2.1), (2.2) moves from an expected initial output (2.3) to an expected final output (2.4) and minimizes the expectation of (2.7). It has been noted in Section 1 that the control based upon $\underline{v} = \underline{0}$ is the same as the control of this problem. Therefore, by the techniques of Sections 3.1 and 3.2, or by optimization of the noise-free problem,

$$2R\underline{u} = B^T \underline{\lambda} , \quad (3.54)$$

$$2H^T P H \dot{\underline{z}} = \dot{\underline{\lambda}} + A^T \underline{\lambda} . \quad (3.55)$$

In expressions (3.20) and (3.21), define a new $\underline{\lambda}_1 = \underline{\lambda}/\mu$, and pass to the limit as $\mu \rightarrow \infty$, assuming $\underline{\lambda}_1$ to exist. The resulting expressions are the same as (3.54) and (3.55), except for the subscript on $\underline{\lambda}_1$.

Solution to the problem of this section is therefore a limiting case of the problem of this paper, when the limit exists. In particular, the analytical results of Section 4, as $\mu \rightarrow \infty$, agree with a direct analysis of (3.54) and (3.55); moreover, the numerical results of Section 3.3.1, for large μ , are in close agreement with those of a corresponding example selected from Section 3.4.

The minimal-variance problem of Section 2 therefore leads to a set of linear, integro-differential, necessary equations. In simple cases, the solution is straightforward; and the numerical results of Section 3.3.1 support (3.29) insofar as the relation of sufficiency conditions and $\mu > 0$ is concerned. Preliminary examples exhibit $4n$ characteristic functions; this property is generalized in the next section to a larger class of problems having only one control variable. Numerical results indicate a reduction in index variance for an increase in index expectation, and vice-versa. Such data is in accordance with the concept of a density function. Finally, the quantities $Q(t,\tau)$, $\underline{q}(t,\tau)$, $q_1(t,\tau)$, and $q_2(t,\tau)$ relate simply to system parameters in these examples. This does not happen in general; and so a general matrix theory of solution is not straightforward.

4. SINGLE-INPUT CONTROL

Suppose there is a single control variable and a single disturbance process, that both enter the plant equations in like manner ($B = C$), and that the output matrix $H = I$. Further, let the v process be white and stationary. Then the necessary conditions (3.27) and (3.28) are equivalent. Wonham and Johnson [19] have found a canonical form for this case. It is shown in this section that there are at most $4n$ characteristic functions associated with the optimal trajectory z and each of its first $n-1$ derivatives, and therefore with the optimal control. It may happen that there are fewer than $4n$ such functions; an example of such a control appears in Section 4.2. It is convenient to assume that v is gaussian.

4.1 Structure of Solution

For a plant ($H = I$) which satisfies (3.28), therefore, there is no loss of generality in a study of the case below:

$$A = \begin{bmatrix} 0 & 1 & 0 & \cdot & \cdot & \cdot & \cdot & 0 \\ 0 & 0 & 1 & 0 & \cdot & \cdot & \cdot & 0 \\ \vdots & & & & & & & \vdots \\ 0 & \cdot & \cdot & \cdot & \cdot & \cdot & 0 & 1 \\ -a_0 & -a_1 & \cdot & \cdot & \cdot & \cdot & \cdot & -a_{n-1} \end{bmatrix} ; \quad (4.1)$$

$$B = \underline{b} = \begin{bmatrix} 0 \\ 0 \\ \vdots \\ 0 \\ 1 \end{bmatrix} ; \quad C = \underline{c} = \begin{bmatrix} 0 \\ 0 \\ \vdots \\ 0 \\ 1 \end{bmatrix} . \quad (4.2)$$

The plant has a characteristic equation

$$\sum_{i=0}^n a_i s^i = 0, \quad a_n = 1; \quad (4.3)$$

and the \underline{z} vector has components

$$z_i = z^{(i-1)}. \quad (4.4)$$

The Laplace transform of the state transition matrix is

$$\phi_{ij}(s) = \frac{\sum_{k=j}^n a_k s^{k+i-j-1}}{\sum_{k=0}^n a_k s^k}, \quad j \geq i; \quad (4.5)$$

$$\phi_{ij}(s) = - \frac{\sum_{k=0}^{j-1} a_k s^{k+i-j-1}}{\sum_{k=0}^n a_k s^k}, \quad j < i. \quad (4.6)$$

If (4.3) has distinct roots $\{s_i\}$, $i = 1, 2, \dots, n$, then

$$\varphi_{ij}(t) = \sum_{m=1}^n \left[\frac{\sum_{k=j}^n a_k s_m^{k+i-j-1}}{\Delta(m)} \right] e^{s_m t}, \quad j \geq i; \quad (4.7)$$

$$\varphi_{ij}(t) = - \sum_{m=1}^n \left[\frac{\sum_{k=0}^{j-1} a_k s_m^{k+i-j-1}}{\Delta(m)} \right] e^{s_m t}, \quad j < i, \quad (4.8)$$

where

$$\Delta(m) = \sum_{k=1}^n k a_k s_m^{k-1}. \quad (4.9)$$

Appropriate changes in (4.7) and (4.8) occur if the s_i are not all distinct. If P is diagonal, having elements P_{ii} , then (3.14), (4.2), and (4.7) give

$$Q_{ij}(t, \tau) = \sigma_v^2 P_{ii} P_{jj} \sum_{m=1}^n \sum_{p=1}^n \left[\frac{s_m^{j-1} s_p^{i-1}}{\Delta(m) \Delta(p) (s_m + s_p)} \right. \\ \left. - \frac{e^{-(s_m + s_p)(t, \tau)}}{s_m + s_p} \right] e^{s_m t} e^{s_p \tau} . \quad (4.10)$$

The first necessary equation (3.20) gives

$$2\mu r u = \lambda_n , \quad (4.11)$$

which combines with the plant equations

$$\sum_{i=0}^n a_i z^{(i)} = u \quad (4.12)$$

to give

$$\lambda_n = 2\mu r \sum_{i=0}^n a_i z^{(i)} . \quad (4.13)$$

Note that the control weighting matrix R is a scalar r . Since v is gaussian, the second necessary equation (3.21) becomes

$$8 \int_0^T Q(t, \tau) \underline{z}(\tau) d\tau + 2\mu P \underline{z} = \dot{\underline{\lambda}} + A^T \underline{\lambda} . \quad (4.14)$$

A systematic elimination of Lagrange variables is possible; and the result is a single integro-differential equation

$$4 \sum_{k=1}^n \sum_{j=1}^n (-1)^{j+1} \frac{\partial^{j-1}}{\partial t^{j-1}} \int_0^T Q_{jk}(t, \tau) z^{(k-1)}(\tau) d\tau \\ + \mu \sum_{j=1}^n (-1)^{j+1} P_{jj} z^{(2j-2)} + \mu r \sum_{j=0}^n \sum_{l=0}^n (-1)^l a_j a_l z^{(j+l)} = 0. \quad (4.15)$$

The extension of the Section 3 methods of differentiation and substitution to (4.15) is difficult. However, it is also possible to proceed as follows. Let

$$z = \sum_{i=1}^{4n} c_i e^{p_i t} ; \quad (4.16)$$

in this case (4.15) becomes

$$\begin{aligned} 0 = & -4\sigma_v^2 \sum_{i=1}^{4n} c_i \sum_{k=1}^n p_i^{k-1} P_{kk} \sum_{j=1}^n (-1)^{j+1} P_{jj} \sum_{p=1}^n \sum_{m=1}^n \left\{ \frac{s_m^{j-1}}{\Delta(m)} \right. \\ & \cdot \frac{s_p^{k-1}}{\Delta(p)} \frac{\partial^{j-1}}{\partial t^{j-1}} \left[\frac{e^{p_i t}}{(p_i - s_m)(p_i + s_p)} \right. \\ & \left. \left. - \frac{e^{s_m t}}{(p_i - s_m)(p_i + s_p)} + \frac{e^{(p_i + s_p)T} (e^{-s_p t} - e^{s_m t})}{(s_m + s_p)(p_i + s_p)} \right] \right\} \\ & + \mu \sum_{i=1}^{4n} c_i \left[\sum_{j=1}^n (-1)^{j+1} p_{ii} p_i^{2j-2} \right] e^{p_i t} \\ & + \mu r \sum_{i=1}^{4n} c_i \left[\sum_{j=0}^n a_j \sum_{l=0}^n (-1)^l a_l p_i^{j+l} \right] e^{p_i t} . \end{aligned} \quad (4.17)$$

The coefficient of $e^{p_i t}$ must vanish identically for each i ; and this occurs if the equation

$$\begin{aligned}
& -4\sigma_v^2 \sum_{k=1}^n p_i^{k-1} p_{kk} \sum_{j=1}^n (-1)^{j+1} p_{jj} p_i^{j-1} \sum_{p=1}^n \sum_{m=1}^n \left[\frac{s_m^{j-1}}{(p_i - s_m) \Delta(m)} \right. \\
& \left. \cdot \frac{s_p^{k-1}}{\Delta(p) (p_i + s_p)} \right] + \mu \sum_{j=1}^n (-1)^{j+1} p_{jj} p_i^{2j-2} \quad (4.18)
\end{aligned}$$

$$+ \mu r \sum_{j=0}^n a_j \left[\sum_{l=0}^n (-1)^l a_l p_i^{j+l} \right] = 0$$

is satisfied by (at most) $4n$ distinct p_i . This expression is, in general, of degree $2n$ in p_i^2 . The coefficients of terms of the form $e^{\pm s_i t}$ in (4.17) must also vanish identically for solution, giving $2n$ conditions. There are also $2n$ conditions on the c_i to meet boundary specifications on $\underline{z}(0)$ and $\underline{z}(T)$. The set of $4n$ equations in the $4n$ unknown c_i appear below.

$$\sum_{i=1}^{4n} c_i p_i^k e^{p_i T} = z^{(k)}(T), \quad k = 0, 1, \dots, n-1; \quad (4.19)$$

$$\sum_{i=1}^{4n} c_i p_i^k = z^{(k)}(0), \quad k = 0, 1, \dots, n-1; \quad (4.20)$$

$$\begin{aligned}
& \sum_{i=1}^{4n} c_i \sum_{k=1}^n p_i^{k-1} p_{kk} \sum_{j=1}^n (-1)^{j+1} p_{jj} \sum_{p=1}^n \left\{ \frac{s_m^{2j-2} s_p^{k-1}}{\Delta(m) \Delta(p)} \right. \\
& \left. \cdot \left[\frac{1}{(p_i - s_m) (p_i + s_p)} + \frac{e^{(p_i + s_p) T}}{(s_m + s_p) (p_i + s_p)} \right] \right\} = 0, \quad (4.21) \\
& m = 1, 2, \dots, n;
\end{aligned}$$

$$\sum_{i=1}^{4n} c_i \sum_{k=1}^n p_i^{k-1} P_{kk} \sum_{j=1}^n (-1)^{j+1} P_{jj} \sum_{m=1}^n \left\{ \frac{s_m^{j-1} s_p^{j+k-2}}{\Delta(m) \Delta(p)} \right\} \cdot \left. \left[\frac{e^{(p_i+s_p)T}}{(s_m+s_p)(p_i+s_p)} \right] \right\} = 0, \quad p = 1, 2, \dots, n. \quad (4.22)$$

If these equations are solvable for the constants c_i , then the vector \underline{z} is given by (4.16) and (4.4); and the open-loop control by (4.12). A necessary and sufficient condition for unique solution is that these equations determine the $\{c_i\}$, $i = 1, 2, \dots, 4n$, uniquely.

4.2 A Degenerate Example

As observed at the beginning of Section 4, it may happen that (4.18) is of degree less than $4n$; that is, the roots of (4.3) have caused cancellations in (4.18). The problem remains determinate, however, since corresponding conditions in the sets (4.21) and (4.22) also vanish. The following example shows a degeneracy of this type.

Let the plant matrix be

$$A = \begin{bmatrix} 0 & 1 \\ -2 & -3 \end{bmatrix}, \quad (4.23)$$

with $s_1 = -1$, $s_2 = -2$. The p_i are roots of

$$p_i^6 - (9 + \frac{1}{r})p_i^4 + [24 + \frac{1}{r}(5 + \frac{4}{\mu})]p_i^2 - [16 + \frac{4}{r}(1 + \frac{1}{\mu})] = 0. \quad (4.24)$$

There are six c_i to determine. Four conditions satisfy specifications on $z(0)$, $z(T)$, $z^{(1)}(0)$, and $z^{(1)}(T)$. The remaining two conditions guarantee that coefficients of e^{+2t} vanish in (4.17). The reader may verify that coefficients of e^{+t} vanish also without further conditions.

For a large class of problems having a single control variable, therefore, the optimal trajectory and control consist in general of at most $4n$ characteristic functions. The uniqueness of the control depends upon the inversion of an equation

$$N \underline{c} = \underline{f} \quad (4.25)$$

for the unknown coefficient vector \underline{c} . Equation (4.25) represents the aggregate of (4.19) through (4.22). A matrix expression for N is not generally available, because of the difficulty in expressing $Q(t, \tau)$ in simple form, or equivalently in performing the integration in (3.25).

5. MULTI-INPUT SYSTEMS

If the control has more than one component, no systematic elimination of Lagrange variables is possible, so that a reduction to scalar equations offers no simplification. Moreover, a matrix analysis depends upon evaluating

$$Q(t, \tau) = PHW(t, \tau)H^T P, \quad (5.1)$$

where

$$W(t, \tau) \triangleq \int_0^{(t, \tau)} \Phi(t-\sigma)CSC^T\Phi^T(t-\sigma)d\sigma. \quad (5.2)$$

Here S is the covariance matrix of the \underline{v} process, which is assumed to be at least wide-sense stationary and white in (5.2). It has not been possible to give a general, algebraic form for (5.2); however, in any specific example the integration is straightforward. It is necessary to discuss the form of $Q(t, \tau)$, rather than its explicit value, in a general formulation.

5.1 A Typical Form for $Q(t, \tau)$

From (5.2), the initial step is to determine the form of the transition matrix. If the eigenvalues of A are distinct, then $\varphi_{ij}(t)$ may be written as

$$\varphi_{ij}(t) = \sum_{k=1}^n c_{kij} e^{s_k t}; \quad (5.3)$$

and the equation

$$W_{ij}(t, \tau) = \int_0^{(t, \tau)} \sum_{k=1}^n \sum_{l=1}^n (CSC^T)_{kl} \varphi_{ik}(t-\sigma) \varphi_{lj}(\tau-\sigma) d\sigma \quad (5.4)$$

becomes

$$W_{ij}(t, \tau) = \sum_{k=1}^n \sum_{\ell=1}^n \sum_{m=1}^n \sum_{p=1}^n \left(\frac{(CSC^T)_{k\ell} c_{mik} c_{plj}}{(s_m + s_p)} \cdot \left[e^{\begin{matrix} s_m t & s_p \tau \\ e^{\begin{matrix} s_p(\tau-t) & t < \tau \\ s_m(t-\tau) & t > \tau \end{matrix} \end{matrix}} \right] \right). \quad (5.5)$$

If the eigenvalues of A are not distinct, then modifications occur in (5.3) and (5.5); such changes do not affect the procedures which follow, except to substitute other linearly independent functions of the type $te^{\beta t}$ or $e^{\alpha t} \cos \beta t$ for the simple exponentials of (5.3).

The use for (5.5) is in the evaluation of the integral

$$\int_0^T Q(t, \tau) e^{P_i \tau} d\tau, \quad (5.6)$$

where $e^{P_i \tau}$ is a scalar function. From the $W(t, \tau)$ matrix, it follows that

$$\left(\int_0^T W(t, \tau) e^{P_i \tau} d\tau \right)_{ij} = \sum_{k=1}^n \sum_{\ell=1}^n \sum_{m=1}^n \sum_{p=1}^n \left(\frac{(CSC^T)_{k\ell} c_{mik} c_{plj}}{(s_m + s_p)} \cdot \left[\left\{ \frac{1}{p_i + s_p} - \frac{1}{p_i - s_m} \right\} e^{p_i t} - \left\{ \frac{e^{(p_i + s_p)T}}{p_i + s_p} \right\} e^{-s_p t} + \left\{ \frac{1}{p_i - s_m} + \frac{e^{(s_p + p_i)T} - 1}{p_i + s_p} \right\} e^{s_m t} \right] \right). \quad (5.7)$$

It is therefore possible to write

$$\begin{aligned}
\int_0^T Q(t, \tau) e^{p_i^T \tau} d\tau &= PHQ_1(p_i) H^T P e^{p_i^T t} \\
&+ \sum_{m=1}^n PHQ_2(p_i, m) H^T P e^{-s_m^T t} \\
&+ \sum_{m=1}^n PHQ_3(p_i, m) H^T P e^{s_m^T t},
\end{aligned} \tag{5.8}$$

where $Q_1(p_i)$, $Q_2(p_i, m)$, and $Q_3(p_i, m)$ must be found for individual examples. The purpose of illustrating the decomposition (5.8) is to show that the fundamental method of solution introduced in the preceding sections is applicable to multi-input systems.

5.2 A Method of Analysis

Let the components of \underline{v} be mutually uncorrelated, and let \underline{v} be gaussian. Then $\underline{q}(t, \tau) = \underline{0}$, and S is diagonal. It is possible to find the eigenvalues of the integro-differential equations (3.10), (3.20), and (3.21) by assuming solutions of the form

$$\underline{z} = \sum_{i=1}^d \underline{c}_{zi} e^{p_i^T t}, \quad \underline{\lambda} = \sum_{i=1}^d \underline{c}_{\lambda i} e^{p_i^T t}, \tag{5.9}$$

where \underline{c}_{zi} and $\underline{c}_{\lambda i}$ are $n \times 1$ constant vectors. Mitra [20] has done related work for integral equations. Then (3.10) and (3.20) combine to give

$$p_i \underline{c}_{zi} = A \underline{c}_{zi} + \frac{1}{2\mu} BR^{-1} B^T \underline{c}_{\lambda i}, \quad i = 1, 2, \dots, d; \tag{5.10}$$

and (3.21) becomes

$$\begin{aligned}
& 8 \sum_{i=1}^d H^T \left(\int_0^T Q(t, \tau) e^{p_i \tau} d\tau \right) H c_{-zi} + \sum_{i=1}^d 2\mu H^T P H c_{-zi} e^{p_i t} \\
& = \sum_{i=1}^d p_i c_{-\lambda i} e^{p_i t} + \sum_{i=1}^d A^T c_{-\lambda i} e^{p_i t} .
\end{aligned} \tag{5.11}$$

From (5.10),

$$c_{-zi} = \frac{1}{2\mu} (p_i I - A)^{-1} B R^{-1} B^T c_{-\lambda i}, \quad i = 1, 2, \dots, d ; \tag{5.12}$$

and a substitution from (5.8) converts (5.11) to

$$\begin{aligned}
& 8 \sum_{i=1}^d D \sum_{k=1}^n [Q_2(p_i, k) e^{-s_k t} + Q_3(p_i, k) e^{s_k t}] D c_{-zi} \\
& + 8 \sum_{i=1}^d D Q_1(p_i) D c_{-zi} e^{p_i t} + 2\mu \sum_{i=1}^d D c_{-zi} e^{p_i t} \\
& = \sum_{i=1}^d (p_i I + A^T) c_{-\lambda i} e^{p_i t} ,
\end{aligned} \tag{5.13}$$

where

$$D = H^T P H . \tag{5.14}$$

If the substitution (5.9) is valid, then the functions $e^{\pm s_k t}$ and $e^{p_i t}$ are linearly independent; and the conditions for solution are that

$$\sum_{i=1}^d D Q_2(p_i, k) D c_{-zi} = \underline{0}, \quad k = 1, 2, \dots, n ; \tag{5.15}$$

$$\sum_{i=1}^d D Q_3(p_i, k) D c_{-zi} = \underline{0}, \quad k = 1, 2, \dots, n ; \tag{5.16}$$

and

$$8DQ_1(p_i)Dc_{-zi} + 2\mu Dc_{-zi} = (p_i I + A^T)c_{\lambda i}, \quad i = 1, 2, \dots, d. \quad (5.17)$$

Then (5.12) combines with (5.17) to give an equation of the form

$$T(p_i)c_{\lambda i} = \underline{0}, \quad (5.18)$$

where

$$T(p_i) = \left[D + \frac{4}{\mu} DQ_1(p_i)D \right] (p_i I - A)^{-1} B R^{-1} B^T - (p_i I - A^T). \quad (5.19)$$

The condition for non-trivial solution is

$$\det T(p_i) = 0, \quad (5.20)$$

which gives an algebraic equation in the eigenvalues $\{p_i\}$, $i = 1, 2, \dots, d$.

Examination of (5.18), (5.20), and (5.12) indicates that $T(p_i)$ has at most a rank of $n-1$, $c_{\lambda i}$ has at least one component which can be freely chosen for each i , and thus that there are at least d constants to be determined. From the boundary conditions (3.11), and the transversality conditions arising from (3.19), $2n$ equations in the unknown constants are obtained. Since D , defined in (5.14), has at most rank r , then (5.15) and (5.16) contribute at most $2nr$ independent conditions. It follows that

$$d_{\max} = 2n(r+1), \quad (5.21)$$

where r is the dimension of \underline{y} , the system output. Since only $2n$ conditions of the form (5.15) and (5.16) occurred in Section 4, d_{\max} was simply $4n$ for that case. It is clear that (5.21) is an equation in d_{\max} , because r appears therein as a maximum estimate.

5.3 Second-Order Example

Suppose that the system is described by the following matrices:

$$H = [1 \quad 1] ; A = \begin{bmatrix} -3 & 0 \\ 0 & -2 \end{bmatrix} ; B = \begin{bmatrix} 1 \\ 1 \end{bmatrix} ; \quad (5.22)$$

$$P = 1 ; C = R = S = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} .$$

This system satisfies (3.28); and

$$\varphi(t) = \begin{bmatrix} e^{-3t} & 0 \\ 0 & e^{-2t} \end{bmatrix} ; \quad (5.23)$$

$$CSC^T = I ; D = H^T P H = \begin{bmatrix} 1 & 1 \\ 1 & 1 \end{bmatrix} .$$

Under these conditions,

$$W(t, \tau) = \begin{bmatrix} \frac{1}{6} [e^{-3|t-\tau|} - e^{-3(t+\tau)}] & 0 \\ 0 & \frac{1}{4} [e^{-2|t-\tau|} - e^{-2(t+\tau)}] \end{bmatrix} ; \quad (5.24)$$

$$Q_1(p_i) = \begin{bmatrix} \frac{1}{9-p_i^2} & 0 \\ 0 & \frac{1}{4-p_i^2} \end{bmatrix} ; \quad (5.25)$$

$$Q_2(p_i, 1) = \begin{bmatrix} \frac{e^{(p_i-3)T}}{6(p_i-3)} & 0 \\ 0 & 0 \end{bmatrix} ; \quad (5.26)$$

$$Q_2(p_i, 2) = \begin{bmatrix} 0 & 0 \\ 0 & \frac{(p_i-2)T}{4(p_i-2)} \end{bmatrix} ; \quad (5.27)$$

$$Q_3(p_i, 1) = \begin{bmatrix} \frac{1}{6} \left[\frac{1-e^{(p_i-3)T}}{p_i-3} - \frac{1}{p_i+3} \right] & 0 \\ 0 & 0 \end{bmatrix} ; \quad (5.28)$$

$$Q_3(p_i, 2) = \begin{bmatrix} 0 & 0 \\ 0 & \frac{1}{4} \left[\frac{1-e^{(p_i-2)T}}{p_i-2} - \frac{1}{p_i+2} \right] \end{bmatrix} ; \quad (5.29)$$

$$\left[D + \frac{4}{\mu} DQ_1(p_i)D \right] = \left[1 + \frac{4}{\mu} \left(\frac{1}{9-p_i^2} + \frac{1}{4-p_i^2} \right) \right] \begin{bmatrix} 1 & 1 \\ 1 & 1 \end{bmatrix} ; \quad (5.30)$$

$$(p_i I - A)^{-1} B R^{-1} B^T = \begin{bmatrix} \frac{1}{p_i+3} & 0 \\ 0 & \frac{1}{p_i+2} \end{bmatrix} \begin{bmatrix} 1 & 1 \\ 1 & 1 \end{bmatrix} . \quad (5.31)$$

Finally,

$$T(p_i) = \begin{bmatrix} \alpha + 3 - p_i & \alpha \\ \alpha & \alpha + 2 - p_i \end{bmatrix} , \quad (5.32)$$

where

$$\alpha = \left(\frac{1}{p_i+3} + \frac{1}{p_i+2} \right) \left(1 + \frac{4}{\mu} \left[\frac{1}{9-p_i} + \frac{1}{4-p_i} \right] \right). \quad (5.33)$$

By (5.20),

$$p_i^2 - p_i(2\alpha + 5) + (5\alpha + 6) = 0, \quad (5.34)$$

which is of eighth degree in p_i , and agrees with (5.21). Moreover, (5.32) does not exist for $p_i = \pm 2, \pm 3$, which validates (5.12). Note also that $\alpha = 0$ is not a possible solution, since (5.18) then demands $p_i = 3, 2$, which gives $\alpha = \infty$ in (5.33), a contradiction. Therefore $T(p_i)$ has rank one; and c_{λ_i} has one free component for each i , or a total of 8 free constants.

From (3.11), $z_1(0)$ and $z_1(T)$ relate to $z_2(0)$ and $z_2(T)$, giving two conditions. Two more conditions arise from the requirement that the conjunct of $\delta J'$ vanish (3.19), so that $\lambda_1(0)$ relates to $\lambda_2(0)$, $\lambda_1(T)$ to $\lambda_2(T)$. It remains to examine the conditions (5.15) and (5.16). Inspection of (5.15), for example, shows only one independent equation for each k , when (5.26) or (5.27) is used; a similar statement holds for (5.16). These then are four conditions. So there are $2n(r+1)$ conditions and $2n(r+1)$ undetermined constants.

The method of analysis is therefore possible for general system matrices. The maximum number of eigenvalues is $2n(r+1)$. This estimate should be compared with the maximum number $4n$ which was found for a single-input, state-output control. Because $Q(t, \tau)$ has not been expressed in algebraic form, the necessary and sufficient conditions for

a unique solution cannot be given in general, but only in terms of the coefficient matrix N in (4.25). Of course N must be determined for each example.

6. CONCLUSIONS

6.1 Summary

The problem of minimizing the variance of a quadratic performance index, in the presence of control noise the properties of which are known a priori, has been studied for linear, constant systems with open-loop control. The method extends in principle to time-varying cases. The necessary (Euler) equations are integro-differential and have a kernel matrix $H^T Q(t, \tau) H$ derived from output correlation functions. Unless the noise is gaussian, white, and has mutually uncorrelated components, these equations contain a forcing vector $H^T \int_0^T \underline{q}(t, \tau) d\tau$ which depends upon the third-moment properties of the disturbance process. For solution, it is necessary that the noise-free system be totally output-controllable with respect to the control matrix B ; if H^{-1} exists, total state controllability with respect to a control matrix CM , where the noise covariance matrix is written as MM^T , guarantees existence of $Q^{-1}(t, \tau)$. Necessary and sufficient conditions for a unique solution may be found for specific examples; a general condition is not possible since an algebraic expression for $Q(t, \tau)$ has not been found.

The expected value of the performance index is constrained to be a positive number K_0 , which cannot be less than the optimal mean index when variance is free. This latter problem, and also the impulsive results when K_0 is free, are limiting cases of this study. The solution has at most $2n(r+1)$ eigenvalues, where n and r are the dimensions of state and output vectors respectively. If state and output are identical, and if there is

a single control variable, this estimate can be refined to $4n$ eigenvalues.

6.2 Problems for Further Study

An open-loop, optimal feedback implementation of these results is possible. This approach modifies the control vector in accordance with observed future states of the system, and so leads to an improvement in performance. The more general problem of minimal-variance, closed-loop control, when inputs are bounded and disturbance statistics are initially unknown, remains unsolved. For the problem of this paper, the question of necessary and sufficient conditions for a unique solution is as yet unresolved. Related to this question is a method for writing $Q(t, \tau)$ in algebraic form, and a procedure for finding the equivalent differential equations for the integro-differential (Euler) set. Further research on the effects of non-gaussian noises and the solution in terms of the forcing function $q(t, \tau)$ will also lead to additional insights.

BIBLIOGRAPHY

1. L. S. Pontryagin, et. al., "The Mathematical Theory of Optimal Processes," Interscience Publishers, New York, N.Y.; 1962.
2. H. J. Kushner, "On the Stochastic Maximum Principle: Fixed Time of Control," Research Institute for Advanced Studies, Baltimore, Md., Tech. Rept. No. 63-24; March, 1964.
3. W. Kipiniak, "Dynamic Optimization and Control: A Variational Approach," The M.I.T. Press and John Wiley & Sons, Inc., New York, N.Y.; 1961.
4. R. E. Bellman, "Adaptive Control Processes: A Guided Tour," Princeton University Press, Princeton, N. J.; 1961.
5. R. E. Bellman and S. E. Dreyfus, "Applied Dynamic Programming," Princeton University Press, Princeton, N. J.; 1962.
6. P. Dorato, R. F. Drenick, and L. Shaw, "Optimal Stochastic Control Theory: A Short Course," Dept. of Electrical Engineering, Polytechnic Institute of Brooklyn, Brooklyn, N. Y.; January, 1964.
7. R. F. Drenick and L. Shaw, "Optimal Control of Linear Plants with Random Parameters," IEEE Trans. on Automatic Control, Vol. AC-9, pp. 236-244; July, 1964.
8. S. E. Dreyfus, "Some Types of Optimal Control of Stochastic Systems," SIAM J. on Control, Vol. 2, No. 1, pp. 120-134; 1964.
9. H. J. Kushner, "Near Optimal Control in the Presence of Small Stochastic Perturbations," Research Institute for Advanced Studies, Baltimore, Md., Tech. Rept. No. 63-22; April, 1964.
10. R. J. Orford, "Optimal Stochastic Control Systems," J. of Mathematical Analysis and Applications, Vol. 6, pp. 419-429; June, 1963.
11. R. E. Kalman, "New Methods in Wiener Filtering Theory," Proc. First Symposium on Engineering Applications of Random Function Theory and Probability, John Wiley & Sons, Inc., New York, N. Y.; 1963.
12. R. Sivan, "The Necessary and Sufficient Conditions for the Optimal Controller to be Linear," Proc. 1964 Joint Automatic Control Conference, Stanford, Calif., pp. 297-304; June, 1964.
13. W. B. Davenport and W. L. Root, "An Introduction to the Theory of Random Signals and Noise," McGraw-Hill Book Co., Inc., New York, N. Y.; 1958.

14. J. L. Doob, "Stochastic Processes," John Wiley & Sons, Inc., New York, N. Y.; 1953.
15. N. Wiener, "Extrapolation, Interpolation, and Smoothing of Stationary Time Series," The M.I.T. Press and John Wiley & Sons, Inc., New York, N. Y.; 1950.
16. E. A. Coddington and N. Levinson, "Theory of Ordinary Differential Equations," McGraw-Hill Book Co., Inc., New York, N. Y.; 1955.
17. E. Kreindler and P. Sarachik, "On the Concepts of Controllability and Observability of Linear Systems," IEEE Trans. on Automatic Control, Vol. AC-9, pp. 129-136; April, 1964.
18. G. A. Bliss, "Lectures on the Calculus of Variations," University of Chicago Press, Chicago, Ill.; 1946.
19. C. D. Johnson and W. M. Wonham, "A Note on the Transformation to Canonical (Phase-Variable) Form," IEEE Trans. on Automatic Control, Vol. AC-9, pp. 312-313; July, 1964.
20. R. Mitra, "On the Solution of a Class of Wiener-Hopf Integral Equations in Finite and Infinite Ranges," Antenna Laboratory, University of Illinois, Urbana, Ill., Tech. Rept. No. 37; 1959.

VITA

Michael Kent Sain was born in Saint Louis, Missouri on March 22, 1937. He received the degrees of Bachelor of Science (June, 1959) and Master of Science (June, 1962) in Electrical Engineering from Saint Louis University. He was a teaching fellow at Saint Louis University from February, 1960 to June, 1962; Paul Galvin Instructor-Fellow at the University of Illinois from September, 1962 to June, 1963; National Science Foundation Summer Fellow in 1963; National Electronics Conference Fellow from September, 1963 to June, 1964; and an Instructor in the Coordinated Science Laboratory at the University of Illinois during the summer of 1964. Currently he is a National Science Foundation Fellow at the University of Illinois. He is a member of Eta Kappa Nu, Pi Mu Epsilon, Alpha Sigma Nu, and Phi Kappa Phi; an associate member of Sigma Xi; and a student member of the Institute of Electrical and Electronics Engineers.

Distribution list as of September 1, 1964

- | | | | | | |
|---|--|----|---|----|---|
| 1 | Director
Air University Library
Maxwell Air Force Base, Alabama
Attn: G-4804 | 1 | Chief of Naval Operations (Code OP-07)
Department of the Navy
Washington, D.C. 20350 | 1 | Commanding General
U.S. Army Electronic Command
Fort Monmouth, New Jersey
Attn: AEGS-RE |
| 1 | Redstone Scientific Information Center
U.S. Army Missile Command
Redstone Arsenal, Alabama | 1 | Commanding Officer
U.S. Army Personnel Research Office
Washington 25, D.C. | 1 | Miss F. Clink
Radio Corporation of America
RCA Laboratories
David Sarnoff Research Center
Princeton, New Jersey |
| 1 | Electronics Research Laboratory
University of California
Berkeley 4, California | 1 | Commanding Officer & Director
Code 140 Library
David M. Taylor Noddi Basin
Washington, D.C. 20360 | 1 | Mr. A.A. Lundstrom
Bell Telephone Laboratories
Room 5E-157
Whippany Road
Whippany, New Jersey |
| 2 | Hughes Aircraft Company
Florence and Teale
Culver City, California
Attn: M.E. Devereux
Technical Document Center | 1 | Chief, Bureau of Ships (Code 686)
Department of the Navy
Washington, D.C. 20360 | 1 | APNOC (DSGSD) Maj. F. Wheeler, Jr.)
Holloman Air Force Base
New Mexico 88530 |
| 3 | Autonics
9150 East Imperial Highway
Downey, California
Attn: Tech. Library, 3041-11 | 1 | Chief, Bureau of Ships (Code 730)
Department of the Navy
Washington, D.C. 20360 | 1 | Commanding General
White Sands Missile Range
New Mexico |
| 1 | Dr. Arnold T. Nordsieck
General Motors Corporation
Defense Research Laboratories
6747 Hollister Avenue
Goleta, California | 1 | Chief, Bureau of Naval Weapons
Technical Library, B11-3
Department of the Navy
Washington, D.C. 20360 | 1 | Microwave Research Institute
Polytechnic Institute of Brooklyn
33 John Street
Brooklyn 1, New York |
| 1 | University of California
Lawrence Radiation Laboratory
P.O. Box 808
Livermore, California | 1 | Director, (Code 5140)
U.S. Naval Research Laboratory
Washington, D.C. 20390 | 1 | Cornell Aeronautical Laboratory, Inc.
4455 Genesee Street
Buffalo 21, New York
Attn: J.F. Desmond, Librarian |
| 1 | Mr. Thomas L. Hartwick
Aerospace Corporation
P.O. Box 95085
Los Angeles 45, California | 1 | Dr. H. Wallace Sinaiko
Institute for Defense Analyses
Research & Engineering Support Division
1664 Connecticut Ave., N.W.
Washington 9, D.C. | 1 | Sperry Gyroscope Company
MARCO Division Library
155 Glen Cove Road
Glen Cove, L.I., New York
Attn: Mrs. Barbara Judd |
| 1 | Professor Zorob Appellian
University of Southern California
University Park
Los Angeles 1, California | 1 | Data Processing Systems Division
National Bureau of Standards
Const. at Van Ness
Room 239, Bldg. 11
Washington 25, D.C.
Attn: A.K. Sallow | 1 | Major William Harris
RADC (RAMT)
Griffiss Air Force Base
New York |
| 1 | Dylvania Electronic Systems-West
Electronic Defense Laboratories
P.O. Box 505
Mountain View, California
Attn: Documents Center | 1 | National Bureau of Standards
Research Information Center &
Advisory Service on Information
Processing
Data Processing Systems Division
Washington 25, D.C. | 1 | Rome Air Development Center
Griffiss Air Force Base
New York
Attn: Documents Library
BAALD |
| 1 | Varian Associates
611 Hansen Way
Palo Alto, California 94303
Attn: Technical Library | 1 | Exchange and Gift Division
The Library of Congress
Washington 25, D.C. | 1 | Library
Light Military Electronics Department
General Electric Company
Armament & Control Products Section
Johnson City, New York |
| 1 | Huston Denlow
Library Supervisor
Jet Propulsion Laboratory
California Institute of Technology
Pasadena, California | 1 | MASA Headquarters
Office of Applications
600 Maryland Avenue, N.W.
Washington 25, D.C.
Attn: Mr. A.M. Greg Andrus
Code FC | 1 | Columbia Radiation Laboratory
Columbia University
538 West 120th Street
New York 27, New York |
| 1 | Professor Nicholas George
California Institute of Technology
Electrical Engineering Department
Pasadena, California | 1 | APQC (PCAFI)
Eglin Air Force Base
Florida | 1 | Mr. Alan Barham
Rome Air Development Center
Griffiss Air Force Base
Rome, New York |
| 1 | Space Technology Labs., Inc.
One Space Park
Redondo Beach, California
Attn: Acquisitions Group
STL Technical Library | 1 | Martin Company
P.O. Box 3437
Orlando, Florida
Attn: Engineering Library
MP-30 | 1 | Dr. E. Howard Holt
Director
Plasma Research Laboratory
Rensselaer Polytechnic Institute
Troy, New York |
| 2 | Commanding Officer and Director
U.S. Navy Electronics Laboratory
San Diego, California 92160
Attn: Code 2800, C.S. Manning | 1 | Commanding Officer
Office of Naval Research, Branch Office
120 North Michigan
Chicago, Illinois 60601 | 3 | Commanding Officer
U.S. Army Research Office (Durham)
Box CM, Duke Station
Durham, North Carolina
Attn: CMD-AA-1P, Mr. Utah |
| 1 | Commanding Officer and Director
U.S. Navy Electronics Laboratory
San Diego, California 92160 | 1 | Laboratories for Applied Sciences
University of Chicago
6200 South Drexel
Chicago, Illinois 60637 | 1 | Battelle-DEFENDER
Battelle Memorial Institute
505 King Avenue
Columbus 1, Ohio |
| 1 | Commanding Officer
Office of Naval Research Branch Office
1500 Geary Street
San Francisco, California 94109 | 1 | Librarian
School of Electrical Engineering
Purdue University
Lafayette, Indiana | 1 | Aeronautical Systems Division
Navigation and Guidance Laboratory
Wright-Patterson Air Force Base
Ohio |
| 1 | The RAND Corporation
1700 Main Street
Santa Monica, California
Attn: Library | 1 | Donald L. Hopley
Department of Electrical Engineering
Iowa State University of Iowa
Iowa City, Iowa | 1 | Aeronautical Systems Division
Directorate of Systems Dynamic Analysis
Wright-Patterson Air Force Base
Ohio |
| 1 | Stanford Electronics Laboratories
Stanford University
Stanford, California
Attn: SEE Documents Librarian | 1 | Commanding Officer
U.S. Army Medical Research Laboratory
Fort Knox, Kentucky | 1 | Commander
Research & Technology Div.
Wright-Patterson Air Force Base
Ohio 45433
Attn: MATT (Mr. Evans) |
| 1 | Dr. L.F. Carter
Chief Scientist Air Force
Room 4E-308, Pentagon
Washington 25, D.C. | 2 | Kests A. Pullen, Jr.
Ballistic Research Laboratories
Aberdeen Proving Ground, Maryland | 1 | Commanding Officer (AD-5)
U.S. Naval Air Development Center
Johnsville, Pennsylvania
Attn: RADC Library |
| 1 | Mr. Robert L. Falk
Associate Director for Research
Research and Technology Division
AFSC
Bolling Air Force Base 25, D.C. | 1 | Director
U.S. Army Human Engineering Laboratories
Aberdeen Proving Ground, Maryland | 2 | Commanding Officer
Frankford Arsenal
Philadelphia 37, Pennsylvania
Attn: SDFPA-1300 |
| 1 | Captain Paul Johnson (USN-Res)
National Aeronautics and Space
Administration
1520 H. Street, N.W.
Washington 25, D.C. | 1 | Commander
Air Force Cambridge Research Laboratories
Laurence G. Hanscom Field
Bedford, Massachusetts | 1 | H.E. Cochran
Oak Ridge National Laboratory
P.O. Box X
Oak Ridge, Tennessee |
| 1 | Major Edwin M. Myers
Headquarters USAF (AFRDR)
Washington 25, D.C. | 1 | Dr. Lloyd Hollingsworth
Director, RSD AFCEC
L.G. Hanscom Field
Bedford, Massachusetts | 1 | U.S. Atomic Energy Commission
Office of Technical Information Extension
P.O. Box 62
Oak Ridge, Tennessee |
| 1 | Dr. James Ward
Office of Deputy Director
(Research and Info)
Department of Defense
Washington 25, D.C. | 1 | Data Sciences Laboratory
Air Force Cambridge Research Lab.
Office of Aerospace Research, USAF
L.G. Hanscom Field
Bedford, Massachusetts
Attn: Lt. Stephen J. Kahne - OSB | 1 | President
U.S. Army Air Defense Board
Fort Bliss, Texas |
| 1 | Dr. Alan T. Weisman
Director, National Science Foundation
Washington 25, D.C. | 1 | Instrumentation Laboratory
Massachusetts Institute of Technology
68 Albany Street
Cambridge 39, Massachusetts
Attn: Library W1-109 | 1 | Director
Human Resources Research Office
The George Washington University
300 North Washington Street
Alexandria, Virginia |
| 1 | Mr. G.D. Watson
Defense Research Member
Canadian Joint Staff
2450 Massachusetts Ave., N.W.
Washington 8, D.C. | 1 | Research Laboratory of Electronics
Massachusetts Institute of Technology
Cambridge 39, Massachusetts
Attn: Document Room 28-307 | 20 | Defense Documentation Center
for Scientific & Technical Information
Cameron Station
Alexandria, Virginia 22314 |
| 1 | Mr. Arthur G. Wimer
Chief Scientist
Air Force Systems Command
Andrews Air Force Base
Washington 25, D.C. | 1 | Dr. Robert Kingston
Lincoln Laboratories
Lexington, Massachusetts | 1 | Commander
U.S. Army Research Office
Highland Building
3045 Columbia Pike
Arlington 4, Virginia |
| 1 | Director Advanced Research
Projects Agency
Washington 25, D.C. | 1 | Lincoln Laboratory
Massachusetts Institute of Technology
P.O. Box 72
Lexington 73, Massachusetts
Attn: Library, A-082 | 1 | U.S. Naval Weapons Laboratory
Computation and Analysis Laboratory
Dahlgren, Virginia
Attn: Mr. Ralph A. Nisemann |
| 1 | Air Force Office of Scientific Research
Directorate of Engineering Sciences
Washington 25, D.C.
Attn: Electronics Division | 1 | Sylvania Electric Products, Inc.
Electronic Systems
Malden Lab., Library
100 First Avenue
Malden 24, Massachusetts | | |
| 1 | Director of Science and Technology
Headquarters, USAF
Washington 25, D.C.
Attn: AFST-84/GV | 1 | Minnesota-Honeywell Regulator Co.
Aeronautical Division
2600 Ridgeway Road
Minneapolis 13, Minnesota
Attn: Dr. D.F. Elwell
Main Station : 605 | | |
| 1 | AFST-8C
Headquarters, USAF
Washington 25, D.C. | 1 | Inspector of Naval Material
Bureau of Ships Technical Representative
1905 West Minnehaha Avenue
St. Paul 4, Minnesota | | |
| 1 | Headquarters, R & T Division
Bolling Air Force Base
Washington 25, D.C.
Attn: RTHR | 20 | Activity Supply Officer, USAERL
Building 2306, Charles Road Area
Fort Monmouth, New Jersey
Fort Accountable Property Officer
Marked: For Inst. for Exploratory Research
Inspect at Destination
Order No. 576-PM-65-91 | | |
| 1 | Headquarters, U.S. Army Materiel Command
Research Division, R & D Directorate
Washington 25, D.C.
Attn: Physics and Electronics Branch
Electronics Section | | | | |
| 1 | Commanding General
U.S. Army Materiel Command
Washington 25, D.C.
Attn: R & D Directorate | | | | |
| 1 | Commanding Officer
Diamond Ordnance Fuse Laboratories
Washington 25, D.C.
Attn: Librarian, Room 211, Bldg. 92 | | | | |
| 1 | Operations Evaluation Group
Chief of Naval Operations (OP-130)
Department of Navy
Washington, D.C. 20350 | | | | |