ON DISCONJUGATE DIFFERENTIAL EQUATIONS

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Introduction. In Wintner's [16] terminology (when n=2), a linear differential equation

$$(0.1) u^{(n)} + p_{n-1}^{(t)} u^{(n-1)} + \cdots + p_0(t) u = 0$$

with continuous coefficients on a t-interval I is said to be disconjugate on I if no solution $(\neq 0)$ has n zeros on I. For n > 2, most known results giving conditions which assure that (0.1) is disconjugate concern perturbations of $u^{(n)} = 0$ either for a fixed finite interval [as, e.g., in the theorem of de la Vallée Poussin (cf. [2, Exercise 5.3(d), p. 346])] or for large t. This paper deals with equations (1.1) which are perturbations of equations with constant coefficients, disconjugate on $-\infty < t < \infty$. As corollaries, we obtain theorems which are refinements of known results concerning perturbations of $u^{(n)} = 0$, but we do not obtain the "best" constants occurring in some of these results (n=2).

The proofs depend on the technique introduced in [5] for discussing asymptotic behavior of solutions of perturbed linear systems with constant coefficients (cf. [2, Chapter X]). This technique is based on suitable changes of variables and arguments which have been subsumed by general theorems of Ważewski [15], (cf. [2, pp. 278–283]). §§1 and 2 use the simple Lemma 4.2, [2, p. 285]; §3 requires a generalization given as Theorem (*) in an Appendix below.

In addition to arguments from the theory of asymptotic integration, the proofs use a theorem of Pólya [14] characterizing equations (0.1) disconjugate on an open interval I in terms of Wronskians of subsets of solutions of (0.1); for a generalization, see [2, pp. 51-54] (also obtained in [6]). Theorems I_{**} and IV of Pólya [14] show that no solution ($\neq 0$) of (0.1) on an open interval I has n distinct zeros if and only if no solution ($\neq 0$) has n zeros counting multiplicities; cf. also [13]. (A generalization of this last fact, when the linear family of solutions of (0.1) is replaced by an arbitrary (not necessarily linear) interpolating family of functions, is given in [1].)

In §4, it is observed that the results of the previous sections, together with theorems and methods of Lasota and Opial [8], give criteria for the existence of solutions of certain nonlinear boundary value problems.

Received by the editors June 27, 1967.

⁽¹⁾ Research partially sponsored by the Air Force Office of Scientific Research, Office of Aerospace Research, United States Air Force, under AFOSR Contract No. AF 49(638)-1382.

1. Distinct real roots. We shall first consider linear differential equations

$$(1.1) u^{(n)} + a_{n-1}u^{(n-1)} + \cdots + a_0u = q_{n-1}(t)u^{(n-1)} + \cdots + q_0(t)u,$$

which are perturbations of equations

$$u^{(n)} + a_{n-1}u^{(n-1)} + \cdots + a_0u = 0$$

with constant coefficients, having real distinct characteristic values.

Assumption (A₁). Let a_0, \ldots, a_{n-1} be n real numbers such that the polynomial

$$\lambda^{n} + a_{n-1}\lambda^{n-1} + \cdots + a_{0} = 0$$

has real distinct roots $\lambda_1 < \lambda_2 < \cdots < \lambda_n$ and let

$$(1.3) c = \min(\lambda_2 - \lambda_1, \ldots, \lambda_n - \lambda_{n-1}) > 0.$$

THEOREM 1.1. Let Assumption (A_1) hold. Then there exists a positive number η_0 , depending on $\lambda_1, \ldots, \lambda_n$ but independent of T, with the following properties: If $q_0(t), \ldots, q_{n-1}(t)$ are continuous functions on $0 \le t < T (\le \infty)$ [or on $0 \le t \le T (< \infty)$] such that

(1.4)
$$q(t) = \left(\sum_{j=0}^{n-1} |q_j(t)|^2\right)^{1/2} \ge 0$$

satisfies

$$\int_0^T q(s) \, ds < \eta_0 \quad or \quad q(t) < c\eta_0$$

or, more generally,

Then (1.1) is disconjugate on $0 \le t < T (\le \infty)$ [or on $0 \le t \le T (< \infty)$].

Proof. Let $y = (u, u', ..., u^{(n-1)})$ and write (1.1) as a linear first order system

(1.6)
$$y' = (A + Q(t))y,$$

where $A = (a_{jk})$ is the constant $n \times n$ matrix with the first n-1 rows given by $a_{jk} = 1$ or $a_{jk} = 0$ for $1 \le j \le n-1$, k = j+1 or $k \ne j+1$ and the last row is $(-a_0, \ldots, -a_{n-1})$, and Q(t) is the matrix with 0 in the first n-1 rows and the last row is $(q_0(t), \ldots, q_{n-1}(t))$.

Let $\Lambda = \Lambda_n = (\lambda_{jk})$ be the constant $n \times n$ matrix with $\lambda_{jk} = \lambda_k^{j-1}$ for $j, k = 1, \ldots, n$. Since the kth column $(1, \lambda_k, \ldots, \lambda_k^{n-1})$ of Λ is an eigenvector of A belonging to the eigenvalue λ_k , we have

(1.7)
$$\Lambda^{-1}A\Lambda = J \equiv \operatorname{diag}(\lambda_1, \ldots, \lambda_n).$$

We make the change of variables

(1.8)
$$y = \Lambda z$$
, i.e., $y_j = \sum_{m=1}^{n} \lambda_m^{j-1} z_m$,

in (1.6) to obtain

$$(1.9) z' = (J+G(t))z,$$

where

(1.10)
$$G(t) = (g_{ik}(t)) = \Lambda^{-1}Q(t)\Lambda.$$

It is clear that $g_{jk}(t)$ is a linear combination of $q_0(t), \ldots, q_{n-1}(t)$ with coefficients depending on j, k and $\lambda_1, \ldots, \lambda_n$. (The explicit formulas for $g_{jk}(t)$ will not be needed below, but a simple calculation shows that

(1.11)
$$g_{jk}(t) = \sum_{m=0}^{n-1} q_m(t) \lambda_k^m / \pi_j,$$

where π_i is defined by

$$\pi_j = (-1)^n \prod_{h>j} (\lambda_h - \lambda_j);$$

cf., e.g., [2, pp. 318-319].) Thus there exists a constant $M_0 = M_0$ $(\lambda_1, \ldots, \lambda_n) > 0$ with the property that

(1.12)
$$\left(\sum_{k=1}^{n} |g_{jk}(t)|^2\right)^{1/2} \leq M_0 q(t) \quad \text{for } j = 1, \ldots, n.$$

(In fact, by (1.11), M_0 can be chosen to be $\|\Lambda\|/\pi$, where $\|\Lambda\|$ is the bound of the matrix Λ as an operator from R^n to R^n and $\pi = \min(|\pi_1|, \ldots, |\pi_n|)$.)

Define the functions

$$\sigma(t) = M_0 \int_0^t q(s)e^{-c(t-s)} ds, \qquad \tau(t) = M_0 \int_t^T q(s)e^{-c(s-t)} ds$$

and assume that the bound η in (1.5) is so small that

$$(1.13) 7M_0\eta < 1.$$

Then, for each k, [2, Lemma 4.2, p. 285] with $\varepsilon = 0$ and $\psi(t) = M_0 q(t)$ and its proof imply that (1.9) has a solution $z = (z_{k1}(t), \ldots, z_{kn}(t)) \neq 0$ satisfying, for $0 \leq t < T$,

$$|z_{kj}(t)| \le 7\sigma(t)z_{kk}(t) \quad \text{for } 1 \le j < k,$$

$$|z_{k'}(t)| \le 7\sigma(t)z_{kk}(t) \quad \text{for } k < j \le n;$$

cf. the Appendix below for a generalization. (Actually, the conditions of Lemma 4.2 in [2, p. 285] and its proof require that $\psi(t) = M_0 q(t) \ge 0$ be positive, but it is clear that the validity of Lemma 4.2 with $\psi(t) > 0$ implies its validity for $\psi(t) \ge 0$.)

By (1.8), the corresponding solution $y = (y_{k1}(t), \dots, y_{kn}(t)) = \Lambda z$ of (1.6) satisfies

$$|y_{kj}(t) - z_{kk}(t)\lambda_k^{j-1}| \le 7M_0 \eta z_{kk}(t) \sum_{m \neq k} |\lambda_m|^{j-1}$$

for $0 \le t < T$. The solution y(t) of (1.6) corresponds to a solution $u = u_k(t)$ of (1.1) with the properties that if $e_{k,l}(t)$ is defined by

$$e_{ki}(t)z_{kk}(t) = u_k^{(j-1)}(t) - \lambda_k^{j-1}z_{kk}(t),$$

then

$$|e_{kj}(t)| \le C\eta$$
, where $C = 7M_0 \max_{j} \sum_{m=1}^{n} |\lambda_m|^{j-1}$.

Let $1 \le m \le n$ and let $W_m(t)$ be the Wronskian determinant of the m solutions $u_1(t), \ldots, u_m(t)$ of (1.1). Then

$$W_m(t) = z_{11}(t) \cdot \cdot \cdot z_{mm}(t) \det \left(\lambda_k^{j-1} + e_{kj} \right)$$

and, consequently, there exists a constant $K = K(\lambda_1, \ldots, \lambda_n)$ with the property that

$$W_n(t) = z_{11}(t) \cdots z_{mm}(t) (\det \Lambda_m + \delta), \qquad |\delta| \leq K \eta (1 + K \eta)^{n-1}$$

and det $\Lambda_m = \prod (\lambda_k - \lambda_j)$, $1 \le j < k \le m$. Hence, there exists an $\eta_0 > 0$ with the property that

$$W_m(t) \neq 0$$
 for $0 \leq t < T, m = 1, \ldots, n$ if $0 \leq \eta < \eta_0$.

It follows from a theorem of Pólya [12, p. 317] that (1.1) is disconjugate on 0 < t < T.

In order to see that (1.1) is disconjugate on [0, T), that is, on every interval [0, b), 0 < b < T, it is only necessary to extend the definition of $q_k(t)$ to an interval $[-\epsilon, T)$, for a suitably small $\epsilon = \epsilon(b) > 0$, and apply the statement already proved with (0, T) replaced by $(-\epsilon, b + \epsilon)$.

COROLLARY 1.1. Let Assumption (A₁) hold. There exists a constant $\eta_0 = \eta_0(\lambda_1, \ldots, \lambda_n) > 0$ with the property that if $q_0(t), \ldots, q_{n-1}(t)$ are continuous for $-\infty < t < \infty$ and q(t) in (1.4) satisfies $q(t) < c\eta_0$ for all t, then (1.1) is disconjugate on $-\infty < t < \infty$.

This is a consequence of Theorem 1.1 which shows that the condition $q(t) < c\eta_0$ implies that (1.1) is disconjugate on every interval [a, b), $-\infty < a < b < \infty$.

COROLLARY 1.2. Let Assumption (A_1) hold. Let $q_0(t), \ldots, q_{n-1}(t)$ be continuous for $0 \le t < \infty$ and q(t) be defined by (1.4). Let $u(t) \ne 0$ be a solution of (1.1) and N = N(T) the number of zeros of u(t), counting multiplicities, on $0 \le t < T$. Then

(1.15)
$$\eta_0 N(T) \leq (n-1) \int_0^T q(t) dt + \eta_0(n-1).$$

Proof. Let the nonnegative zeros (if any) of u(t) be $0 \le t_0 \le t_1 \le \cdots$. Then

$$\int_{s}^{t} q(\rho) d\rho \ge \eta_0 \quad \text{if } s = t_j, t = t_{j+n-1}.$$

Thus if $t_{N-1} < T \le t_N$,

$$\int_0^T q(\rho) \ d\rho \ge \eta_0[(N-1)/(n-1)] \ge \eta_0\{(N-1)/(n-1)-(n-2)/(n-1)\}.$$

This gives (1.15).

2. Coincident real roots. We now consider the case when Assumption (A_1) does not hold and, in fact, the extreme opposite case, $\lambda_1 = \cdots = \lambda_n$, holds. We suppose that the value of $\lambda_1 = \cdots = \lambda_n$ is 0; cf. the Remark at the end of this section.

THEOREM 2.1. There exists a number $\eta_1 > 0$, independent of T, with the following property: In the differential equation

$$(2.1) u^{(n)} = q_{n-1}(t)u^{(n-1)} + \cdots + q_0(t)u,$$

let $q_0(t), \ldots, q_{n-1}(t)$ be continuous for $0 \le t < T \ (\le \infty)$ [or for $0 \le t \le T \ (< \infty)$] and such that

(2.2)
$$q(t) = \left(\sum_{k=0}^{n-1} |q_k(t)t^{n-k-1}|^2\right)^{1/2} \ge 0$$

satisfies

(2.3)
$$t^{-1} \int_0^t sq(s) \, ds, \qquad t \int_t^T s^{-1}q(s) \, ds \leq \eta < \eta_1$$

(in particular, let

(2.4)
$$\int_0^T q(t) dt < \eta_1 \quad or \quad tq(t) < \eta_1$$

hold), then (2.1) is disconjugate on $0 \le t < T (\le \infty)$ [or on $0 \le t \le T (< \infty)$].

If n=2 and $q_1(t)\equiv 0$, then the first criterion in (2.4) is known with $\eta_1=1$ (cf. [4, Theorem 5.1, p. 345] with m(t)=t-a) and the second is an analogue of A. Kneser's criterion $t^2q_0(t)\leq 1/4$. If $n\geq 2$, $q_1(t)\equiv \cdots \equiv q_{n-1}(t)\equiv 0$ and $q_0(t)$ is of constant sign and monotone, then the first criterion in (2.4) has been proved by Nehari [12]. (For a related result of de la Vallée Poussin and generalizations, see Levin [9], [10], Nehari [11], and Hukuhara [7].) According to a theorem of Dunkel (cf. [2, Theorem 17.1, p. 315]), the convergence of the integral in (2.4) for $T=\infty$ implies asymptotic formulae for the solutions of (2.1) for large t. But these formulae do not imply a disconjugacy criterion for (2.1) on $t\geq 0$.

Proof. Under the change of independent variables

$$(2.5) s = -\log t or t = e^{-s},$$

the interval 0 < t < T is changed to $-\log T < s < \infty$. Note that $u' = -e^s Du$, where D = d/ds. Thus

$$u^{(k)} = (-t)^{-k}D(D+1)\cdots(D+k-1)u.$$

Write this relation as

$$u^{(k)} = (-t)^{-k} \sum_{j=1}^{k} \beta_{kj} D^{j}$$
, where $\beta_{kk} = 1$,

so that

(2.6)
$$\sum_{j=1}^{k} \beta_{kj} \lambda^{j} = \lambda(\lambda+1) \cdots (\lambda+k-1).$$

Thus (2.1) becomes

(2.7)
$$D^{n}u + \sum_{j=1}^{n-1} \beta_{nj}D^{j}u = Q_{n-1}(s)D^{n-1}u + \cdots + Q_{0}(s)u,$$

where $t = e^{-s}$, $Q_0(s) = (-t)^n q_0(t)$, and

(2.8)
$$Q_{j}(s) = \sum_{k=j}^{n-1} (-t)^{n-k} q_{k}(t) \beta_{kj} \quad \text{for } j = 1, \dots, n-1.$$

Hence, there is a constant C = C(n) such that

(2.9)
$$Q(s) = \left(\sum_{j=0}^{n-1} |Q_j(s)|^2\right)^{1/2} \leq Cq(t)t,$$

where q(t) is defined by (2.2).

It is clear from (2.6) that the roots of the characteristic polynomial belonging to the left side of (2.7) are $0, -1, \ldots, -(n-1)$. In view of (2.9), Theorem 2.1 follows from Theorem 1.1 after an obvious change of integration variables.

Theorem 2.1 will be combined with the method of [3] for the case n=2, $q_1(t)\equiv 0$ (cf. [2, pp. 346-347]) to obtain an analogue of Corollary 1.2.

COROLLARY 2.1. There exist positive constants M_0, \ldots, M_{n-1} with the following property: Let $q_0(t), \ldots, q_{n-1}(t)$ be continuous for $t \ge 0$, $u(t) \ne 0$ a solution of (2.1), and N = N(T) the number of zeros of u(t) on $0 \le t < T$ ($< \infty$), counting multiplicities. Then $N \le n-1$ or N satisfies the inequality

$$(2.10) \qquad \sum_{k=0}^{n-1} M_k \{ N/(n-1) - 1 \}^{k-n} \left\{ T^{n-1} \int_0^T |q_k(t)|^{n/(n-k)} dt \right\}^{1-k/n} > 1.$$

Since k-n<0, (2.10) can be used to estimate N=N(T) from above. An admissible choice for M_0, \ldots, M_{n-1} is given by $M_k=2^{k-n}, k=0, \ldots, n-1$. This follows from the proof of Corollary 2.1 and a result of Nehari [11]. (Another choice of the constants may be given in the paper by Hukuhara [7] which is not available to me; the pertinent (last) result stated in the review in the Mathematical Reviews does

not seem correct.) Still another choice is, e.g., $M_0 = (n-2)^{n-2}/(n-1)^{n-1}(n-1)!$ and $M_k = 1/(n-k-1)!$ for $k=1,\ldots,n-1$; cf. the argument in [2, Exercise 5.3, p. 570]. See also Levin [9], [10].

Proof. Let N > n-1 and $(0 \le)$ $t_0 \le t_1 \le \cdots$ be the nonnegative zeros of u(t); so that $t_{N-1} < T \le t_N$. By Theorem 2.1, there exist positive constants M_0, \ldots, M_{n-1} such that

$$\sum_{k=0}^{n-1} M_k(t-s)^{-k} \int_s^t |q_k| d\rho > (t-s)^{1-n} \quad \text{if } s = t_j, t = t_{j+n-1};$$

e.g., by (2.2) and (2.4), any set of constants $M_k > 1/\eta_1$ for k = 0, ..., n-1 is admissible. By [11], one can choose $M_k = 2^{-k}$ for k = 0, ..., n-1.

Letting $s = t_i$, $t = t_{i+n-1}$ and adding for j = 0, ..., m gives

(2.11)
$$\sum_{k=0}^{n-1} M_k \sum_{j=0}^{m} (t-s)^{-k} \int_{s}^{t} |q_k| d\rho > \sum_{j=0}^{m} (t-s)^{1-n}.$$

Note that if $1 < \alpha < \infty$,

$$(t-s)^{-k} \int_{s}^{t} |q_{k}| d\rho \leq (t-s)^{1-k-1/\alpha} \left(\int_{s}^{t} |q_{k}|^{\alpha} d\rho \right)^{1/\alpha}$$

Thus, if $1/\alpha + 1/\beta = 1$,

$$\sum_{j=0}^{m} (t-s)^{-k} \int_{s}^{t} |q_{k}| d\rho \leq \left(\sum_{j=0}^{m} (t-s)^{\beta(1-k-1/\alpha)} \right)^{1/\beta} \left(\int_{t_{0}}^{t_{m(n-1)}} |q_{k}|^{\alpha} d\rho \right)^{1/\alpha}.$$

For $1 \le k \le n-1$, choose $\alpha = n/(n-k)$ and $\beta = n/k$, so that $\beta(1-k-1/\alpha) = 1-n$ and the last product is

$$\left(\sum_{i=0}^{m} (t-s)^{1-n}\right)^{k/n} \left(\int_{t_0}^{t_{m(n-1)}} |q_k|^{n/(n-k)} d\rho\right)^{1-k/n}.$$

Thus, by (2.11),

$$\sum_{k=0}^{n-1} M_k \left(\sum_{j=0}^m (t-s)^{1-n} \right)^{-1+k/n} \left(\int_{t_0}^{t_{m(n-1)}} |q_k|^{n/(n-k)} d\rho \right)^{1-k/n} > 1.$$

From the inequality for the harmonic mean,

$$\left[\sum_{t=0}^{m} (t-s)^{-(n-1)}\right]^{-1/(n-1)} \leq (m+1)^{-n/(n-1)} \sum_{t=0}^{m} (t-s).$$

The last sum is $t_{m(n-1)} - t_0 < T$ if $t_{m(n-1)} \le t_{N-1} < T$. Thus we obtain

$$\sum_{k=0}^{n-1} M_k \left[(m+1)^{-n} T^{n-1} \int_0^T |q_k|^{n/(n-k)} ds \right]^{1-k/n} > 1.$$

If we choose $m+1=[(N-1)/(n-1)] \ge (N-1)/(n-1)-(n-2)/(n-1)$, then $m+1 \ge N/(n-1)-1$. This, together with the last inequality, gives (2.10) and proves Corollary 2.1.

REMARK. The analogues of Theorem 2.1 and its corollary in which (2.1) is replaced by

$$(2.12) (d/dt - \lambda)^n u = q_{n-1}(t)u^{(n-1)} + \cdots + q_0(t)u,$$

and (2.2) by

(2.13)
$$q(t) = \left(|q_0(t)t^{n-1}|^2 + \sum_{k=1}^{n-1} |q_k(t)t^{n-k-1}(1+|\lambda t|^k)|^2 \right)^{1/2},$$

are valid (and the corresponding number η_1 does not depend on λ). In fact, the change of dependent variables

$$v = e^{-\lambda t}u$$

transforms (2.12) into an equation of the form

$$v^{(n)} = \rho_{n-1}(t)v^{(n-1)} + \cdots + \rho_0(t)v,$$

where

$$\rho_j(t) = \sum_{k=1}^{n-1} C_{kj} \lambda^{k-j} q_k(t), \text{ where } C_{kj} = k!/j!(k-j)!.$$

Hence, (2.13) satisfies

$$\sum_{k=0}^{n-1} |\rho_k(t)t^{n-k-1}|^2 \le C^2 |q(t)|^2$$

for a suitable constant C, independent of λ .

3. Arbitrary real roots. We now consider the differential equation

$$(3.1) u^{(n)} + a_{n-1}u^{(n-1)} + \cdots + a_0u = q_{n-1}(t)u^{(n-1)} + \cdots + q_0(t)u,$$

when the following holds:

ASSUMPTION (A₂). Let a_0, \ldots, a_{n-1} be real numbers such that the polynomial

(3.2)
$$\lambda^{n} + a_{n-1}\lambda^{n-1} + \cdots + a_{0} = 0$$

has only real roots, say, $\lambda(1), \ldots, \lambda(g)$, with the respective multiplicities $h(1), \ldots, h(g)$; so that $h(j) \ge 1$ and $h(1) + \cdots + h(g) = n$. Let

(3.3)
$$h_* = \max(h(1), \ldots, h(g)).$$

THEOREM 3.1. Let Assumption (A_2) hold. Then there exist numbers $C \ge 1$ and $\eta_2 > 0$, depending on a_0, \ldots, a_{n-1} but independent of T, with the following property: If $q_0(t), \ldots, q_{n-1}(t)$ are continuous on $0 \le t < T$ ($\le \infty$) [or on $0 \le t \le T$ ($< \infty$)] and

(3.4)
$$q(t) = (1+t)^{h_0-1} \sum_{k=0}^{n-1} |q_k(t)|$$

satisfies

(3.5)
$$\int_0^t q(s)[(s+C)/(t+C)] ds, \int_t^T q(s)[(t+C)/(s+C)] ds \leq \eta < \eta_2$$

or, in particular,

then (3.1) is disconjugate on $0 \le t < T (\le \infty)$ [or on $0 \le t \le T (< \infty)$].

The proof depends on a method of Hartman and Wintner [5] for the asymptotic integration of linear systems with nearly constant coefficients. This method involves a change of variables which enables us to obtain rather precise information on the asymptotic behavior of each component of the solution vector.

Proof. We shall deal only with disconjugacy for 0 < t < T. The passage to $0 \le t < T$ can be made as in the proof of Theorem 1.1.

(a) Let $C = C(a_0, \ldots, a_{n-1}) \ge 1$ be a number to be fixed below and introduce the abbreviation

(3.7)
$$s = t + C \ge 1$$
 if $0 \le t < T$.

Assume, that for each j, there exist h(j) solutions $u = u_{j1}(t), \ldots, u_{jh(j)}(t)$ of (3.1) with the properties that $v_{j\kappa}(t) = u_{j\kappa}(t)e^{-\lambda(j)s}$ and its derivatives satisfy

$$|D^{i-1}v_{j\kappa}(t) - w_{\kappa}^{j}(t)s^{\kappa-i}/(\kappa-1)!| \leq \gamma\eta|w_{\kappa}^{j}(t)|s^{\kappa-i} \quad \text{for } i = 1, \ldots, \kappa,$$

$$|D^{i-1}v_{j\kappa}(t)| \leq \gamma\eta|w_{\kappa}^{j}(t)|s^{\kappa-i} \quad \text{for } i = \kappa+1, \ldots, h(j),$$

$$|D^{i-1}v_{j\kappa}(t)| \leq \gamma\eta|w_{\kappa}^{j}(t)|s^{\kappa-h_{\bullet}} \quad \text{for } h(j) < i \leq n,$$

for $0 \le t < T$, where $w_k^j(t)$ is a continuous nonvanishing function; γ denotes a constant, not always the same, depending only on a_0, \ldots, a_{n-1} ; and D = d/dt.

It will now be shown that if $\eta > 0$ is sufficiently small and the *n* solutions $u_{11}, \ldots, u_{1h(1)}, u_{21}, \ldots, u_{gh(g)}$ are numbered as u_1, \ldots, u_n , then the Wronskian determinant

$$W_m(t) = \det(D^{i-1}u_j), \text{ where } i, j = 1, ..., m (\leq n),$$

does not vanish for $0 \le t < T$.

Let $1 \le m \le n$; say, $m = h(1) + \cdots + h(\mu - 1) + \nu$, where $1 \le \nu \le h(\mu)$. Let h'(j) = h(j) or ν according as $j \le \mu - 1$ or $j = \mu$. The entry $D^{i-1}u_{j\kappa}$ in $W_m(t)$ is

(3.9)
$$D^{i-1}u_{j\kappa} = D^{i-1}e^{\lambda(j)s}v_{j\kappa} = e^{\lambda(j)s}\sum_{r=1}^{i}C_{i-1,r-1}\lambda^{i-r}(j)D^{r-1}v_{j\kappa}.$$

If $i \le h(j)$, let $e_{ih\kappa}(t) = 0$. If i > h(j), let

(3.10)
$$e_{ij\kappa}(t) = \sum_{r=h(j)+1}^{i} C_{i-1,r-1} \lambda^{i-r}(j) D^{r-1} v_{j\kappa},$$

so that, by (3.8),

$$(3.11) |e_{ij\kappa}(t)| \leq \gamma |w_{\kappa}^{i}(t)| \eta s^{\kappa - h_{\bullet}} \text{ for } 0 \leq t < T; \kappa = 1, \ldots, h(j); i = 1, \ldots, n.$$

The relations (3.9), (3.10) give

(3.12)
$$D^{i-1}u_{j\kappa} = e^{\lambda(j)s}[f_{ij\kappa}(t) + e_{ij\kappa}(t)],$$

where

$$f_{ij\kappa} = \sum_{r=1}^{\min(i,h(j))} C_{i-1,r-1} \lambda^{i-r}(j) D^{r-1} v_{j\kappa}.$$

The $m \times m$ matrix $(f_{ijk}(t))$, corresponding to $(D^{i-1}u_j)$, $i, j=1, \ldots, m$, can be written as the product $\Lambda_{(m)}W^m(t)$ of two $m \times m$ matrices. The matrix $\Lambda_{(m)}$ is a constant matrix which is the Wronskian determinant at t=0 of the first m functions of $\omega_{11}, \ldots, \omega_{1h(1)}, \omega_{21}, \ldots, \omega_{gh(g)}$, where $\omega_{jk} = e^{\lambda(j)t}t^{k-1}/(k-1)!$. By (3.9),

$$D^{i-1}\omega_{jr}(0) = \sum_{k} C_{i-1,k-1}\lambda^{i-k}(j)[t^{r-k}/(r-k)!]_{t=0},$$

and the sum is over the range $1 \le k \le \min(i, r)$; so that

$$D^{i-1}\omega_{ir}(0) = C_{i-1,r-1}$$
 or 0, according as $r \le i$ or $r > i$.

Choose $W^m(t)$ to be the matrix

$$(3.14) W^m(t) = \operatorname{diag}(W_{(1)}(t), \ldots, W_{(\mu)}(t)),$$

where $W_{(j)}(t)$ is an $h'(j) \times h'(j)$ matrix,

$$W_{(j)}(t) = (D^{i-i}v_{j\kappa})$$
 for $i, \kappa = 1, ..., h'(j)$.

Thus, we have $(f_{ijk}(t)) = \Lambda_{(m)} W^m(t)$.

The matrix $\Lambda_{(m)}$ is nonsingular, since its determinant is the Wronskian determinant of m linearly independent solutions of an mth order linear equation $u^{(m)} + \cdots = 0$ with constant coefficients. Thus, by (3.12), the matrix $(D^{i-1}u_j)$ can be written

$$(D^{i-1}u_{i\kappa}) = \Lambda_{(m)}[W^m(t) + \Lambda_{(m)}^{-1}(e_{ij\kappa})]D_0,$$

where $D_0 = \text{diag}(e^{\lambda(1)s}, \ldots, e^{\lambda(\mu)s})$ in which $e^{\lambda(j)s}$ occurs h'(j) times. Thus, the assertion $W_m(t) \neq 0$ is equivalent to

(3.15)
$$\det W^{(m)}(t) \neq 0, \text{ where } W^{(m)} = W^m(t) + \Lambda_{(m)}^{-1}(e_{ijk}).$$

From (3.11), the absolute values of the elements in the columns of $\Lambda_{m_1}^{-n_1}(e_{ij\kappa})$ corresponding to $u_{j\kappa}$ are majorized by $\gamma|w_{\kappa}^j(t)|\eta s^{\kappa-h_{\star}}$. Divide the corresponding column of $W^{(m)}(t)$ by $s^{\kappa-1}|w_{\kappa}^j(t)|$ and multiply the $[h(1)+\cdots+h(j-1)+i]$ th row by s^{i-1} . In the resulting matrix, the elements which are not in the blocks corresponding to $W_{(1)},\ldots,W_{(\mu)}$ (cf. (3.14)) are majorized by $\gamma\eta s^{i-h_{\star}}$, where $i-h_{\star} \leq h(j)-h_{\star} \leq 0$. Since $s \geq 1$ for $t \geq 0$, $\gamma\eta s^{i-h_{\star}} \leq \gamma\eta$. By the same argument, the elements below the diagonals on the blocks $W_{(1)},\ldots,W_{(\mu)}$ are majorized by $\gamma\eta s^{i-\kappa} \leq \gamma\eta$ for $i \leq \kappa$. The elements on and above the diagonal become $1/(i-\kappa)!$ + error, where $|\text{error}| \leq \gamma\eta$; in particular, the diagonal elements differ from 1 by at most $\gamma\eta$. Thus, it is clear that if η_2 is sufficiently small and $\eta < \eta_2$, then (3.15) holds.

(b) In order to complete the proof of Theorem 3.1, it is necessary to prove the existence of the solutions $u=u_{j\kappa}(t)$ assumed in part (a). Let j be fixed, $1 \le j \le g$. For the sake of definiteness let j=1, $\lambda=\lambda(1)$, h=h(1) and put $\mu(j)=\lambda(j)-\lambda$ for $j=1,\ldots,g$. In particular $\mu(1)=0$. In (3.1), introduce the new dependent variable

$$(3.16) v = e^{-\lambda t}u.$$

Then (3.1) takes the form

$$(3.17) v^{(n)} + b_{n-1}v^{(n-1)} + \cdots + b_nv^{(n)} = \sum_{i=0}^{n-1} \rho_i(t)v^{(i)}.$$

Here b_n, \ldots, b_{n-1} are constants such that the roots of

$$\tau^n + b_{n-1}\tau^{n-1} + \cdots + b_n\tau^n = 0$$

are $0 = \mu(1), \mu(2), \dots, \mu(g)$ with the respective multiplicities $h(1), \dots, h(g)$. Also

(3.18)
$$\rho_i(t) = \sum_{k=i}^{n-1} C_{k,i} \lambda^{k-i} q_k(t), \text{ where } C_{k,i} = k!/i!(k-i)!.$$

Write (3.17) as a first order system

$$(3.19) y' = By + P(t)y$$

for the vector $y = (v, v', ..., v^{(n-1)})$; cf. (1.1) and (1.6).

For a given j, let J_j denote the $h(j) \times h(j)$ Jordan matrix with the diagonal elements $\mu(j)$ and, if h(j) > 1, superdiagonal elements 1. The first h columns of B constitute an $h \times n$ matrix with J_1 in the upper portion and the zero matrix $0_{h \times (n-h)}$ below. Thus, there exists a nonsingular constant matrix of the form

(3.20)
$$\Lambda = \begin{pmatrix} I_h & \cdots \\ 0 & \cdots \end{pmatrix}, \qquad \Lambda^{-1} = \begin{pmatrix} I_h & \cdots \\ 0 & \cdots \end{pmatrix},$$

such that I_h is the unit $h \times h$ matrix and

(3.21)
$$\Lambda^{-1}B\Lambda = J \equiv \operatorname{diag}(J_1, \ldots, J_g).$$

The linear change of variables

$$(3.22) y = \Lambda z$$

sends (3.19) into

$$(3.23) z' = Jz + \Lambda^{-1}P(t)\Lambda z.$$

Write $z = (z^1, \ldots, z^g)$, where $z^j = (z^j_1, \ldots, z^j_{h(g)})$ is a vector of dimension h(g), and (3.23) as

(3.24)
$$z^{j'} = J_j z^j + \sum_{k=1}^g P^{jk}(t) z^k \text{ for } j = 1, \dots, g,$$

where $P^{jk}(t)$ is an $h(j) \times h(k)$ matrix.

Since P(t) is an $n \times n$ matrix with 0 entries except for the last row $(-\rho_0(t), \ldots, -\rho_{n-1}(t))$, it follows from (3.20), (2.23) that the *i*th component of $P^{j1}(t)z^1$ is of the form

$$(3.25) (P^{j1}(t)z^1)_i = \alpha_i^1 \sum_{m=0}^{h-1} \rho_m(t)z_{m+1}^1,$$

where $(-\alpha_1^1, \ldots, -\alpha_{h(g)}^g)$ is the last column of Λ^{-1} . Furthermore the elements of $P^{jk}(t)$ are linear combinations of $\rho_0(t), \ldots, \rho_{n-1}(t)$ with constant coefficients.

Let $C \ge 1$, $\varepsilon > 0$ be constants to be specified below and introduce the abbreviation s = t + C, as in (a) above. (Note that t is still the independent variable.) Make the linear change of variables $z \to w$:

(3.26)
$$z^{1} = (\exp J_{1}s)w^{1} \text{ and } z_{i}^{j} = \varepsilon^{h(j)-i}w_{i}^{j}$$

for $j=2,\ldots,g$, to obtain

(3.27)
$$w^{1'} = \sum_{k=1}^{g} Q^{1k}(t)w^k \text{ and } w^{j'} = J_{j\varepsilon}w^j + \sum_{k=1}^{g} Q^{jk}(t)w^k,$$

where $J_{j\varepsilon}$ is obtained by replacing the superdiagonal elements 1 in J_j by ε . If (3.26) is written as

$$z = Q_0(t)w, Q_0(t) = \text{diag}(\exp J_1 s, D_1),$$

where $D_1 = D_1(\varepsilon)$ is a diagonal matrix, then (3.27) is

$$w' = \operatorname{diag}(0, J_{2\varepsilon}, \ldots, J_{g\varepsilon})w + Q_0^{-1}\Lambda^{-1}P\Lambda Q_0w.$$

From this, we can write (3.27) in detail as

$$(3.28) w_i^{1'} = \left(\sum_{i=1}^h c_{ir} s^{r-i}\right) \sum_{k=1}^h \sum_{m=0}^{k-1} d_{km} \rho_m(t) s^{k-1-m} w_k^1 + \sum_{l=0}^g \sum_{m=1}^{h(1)} \sum_{m=1}^h L_{1i,lmr}(t) s^{r-i} w_m^l,$$

 $i=1,\ldots,h$, and

$$(3.29) w_i^{j'} = \mu_j w_i^j + \varepsilon_{ji} w_{i+1}^j + \sum_{k=1}^h \sum_{m=0}^{k-1} d_{ji,km} \rho_m(t) s^{k-1-m} w_k^1 + \sum_{l=2}^g \sum_{m=1}^{h(1)} L_{ji,lm}(t) w_m^l,$$

for $j=2,\ldots,g$ and $i=1,\ldots,h(j)$, where $c_{ir},d_{km},d_{ji,km}$ are constants; $L_{1i,lmk}(t)$ and $L_{ji,lm}(t)$ are linear functions of $\rho_0(t),\ldots,\rho_{n-1}(t)$ with constant coefficients; and $\varepsilon_{ji}=\varepsilon$ or 0 according as i< h(j) or i=h(j).

Let $1 \le \kappa \le h$. Make the last change of variables $w \to x$:

(3.30)
$$x_i^1 = s^{i-\kappa} w_i^1 \text{ and } x_i^j = s^{h_{\bullet}-\kappa} w_i^j, \quad j > 1.$$

Then (3.28)-(3.29) becomes

$$(3.31) \quad x_i^{1'} = (i - \kappa) s^{-1} x_i^1 + \sum \sum \sum c_{ir} d_{km} \rho_m(t) s^{r-m-1} x_k^i + \sum \sum \sum L_{1i,lmr}(t) s^{r-h_0} x_m^l,$$

(3.32)
$$x_{i}^{j'} = [\mu(j) + (h_{*} - \kappa)s^{-1}]x_{i}^{j} + \varepsilon_{ji}x_{i+1}^{j} + \sum_{j} \sum_{l \neq i, km} \rho_{m}(t)s^{h_{*} - m - 1} + \sum_{j} \sum_{l \neq i, lm} (t)x_{m}^{l},$$

where the ranges for the sums are the same as in (3.28), (3.29), respectively.

Choose $C \ge 1$ so large and $\varepsilon > 0$ so small that

(3.33)
$$[\operatorname{sgn} \mu(j)] J_{is} x^j \cdot x^j - (h_* - \kappa) s^{-1} |x^j|^2 \ge s^{-1} |x^j|^2 \quad \text{for } j > 1.$$

Since $s = t + C \ge 1$ when $t \ge 0$, it is clear that there exists a constant

$$C_1 = C_1(a_0, \ldots, a_{n-1})$$

such that if

(3.34)
$$\psi(t) = C_1 \sum_{m=0}^{n-1} |\rho_m(t)| s^{h_{\bullet}-1},$$

then

$$|\{\cdots\}| \le \psi(t)|x|$$
 if $|x|^2 = \sum_{j=1}^g |x^j|^2 = \sum_{j=1}^g \sum_{k=1}^{h(j)} |x_k^j|^2$

and $\{\cdots\}$ denotes either the sum of the two triple sums in (3.31) or the sum of the two double sums in (3.32). Thus, (3.31)–(3.32) imply that

$$\sum_{i=1}^{\kappa-1} x_i^1 x_i^{1'} \leq -s^{-1} \sum_{i=1}^{\kappa-1} |x_i^1|^2 + \psi(t)|x| \left(\sum_{i=1}^{\kappa-1} |x_i^1|^2 \right)^{1/2},$$

$$x^j \cdot x^{j'} \leq -s^{-1}|x^j|^2 + \psi(t)|x| \cdot |x^j| \quad \text{if } \mu(j) < 0,$$

$$|x_\kappa^1 x_\kappa^{1'}| \leq \psi(t)|x| \cdot |x_\kappa^1|,$$

$$\sum_{i=\kappa+1}^h x_i^1 x_i^{1'} \geq s^{-1} \sum_{i=\kappa+1}^h |x_i^1|^2 - \psi(t)|x| \left(\sum_{i=\kappa+1}^h |x_i^1|^2 \right)^{1/2},$$

$$x^j \cdot x^{j'} \geq s^{-1}|x^j|^2 - \psi(t)|x| \cdot |x^j| \quad \text{if } \mu(j) > 0.$$

We can now apply Theorem (*) of the Appendix below (with y^{κ} corresponding to the 1-dimensional x_{κ}^{1} with $\alpha^{\kappa} = \alpha_{\kappa} = 0$; the vectors $y^{1}, \ldots, y^{\kappa-1}$ correspond to $y^{1} = (x_{1}^{1}, \ldots, x_{\kappa-1}^{1})$ and the x^{j} for which $\mu(j) < 0$ with the respective $\alpha_{1} = \alpha_{2} = \cdots = \alpha_{\kappa-1} = -s^{-1}$; the vectors $y^{\kappa+1}, \ldots, y^{n}$ correspond to $y^{\kappa+1} = (x_{\kappa+1}^{1}, \ldots, x_{n}^{1})$ and the x^{j} for which $\mu(j) > 0$ with the respective $\alpha_{\kappa+1} = \alpha_{\kappa+2} = \cdots = s^{-1}$.

Thus, we have to examine the functions

$$\sigma_1 = \int_0^t \psi(\rho)(\rho+C)/(t+C) \ d\rho, \qquad \sigma_{\kappa+1} = \int_t^T \psi(\rho)(t+C)/(\rho+C) \ d\rho.$$

According to Theorem (*), if σ_1 and $\sigma_{\kappa+1}$ exist and have a sufficiently small bound η , then (3.31)–(3.32) has a solution $x(t) \not\equiv 0$ such that

$$|x_i^1| \le \gamma \eta |x_{\kappa}^1|$$
 for $i \ne \kappa$, $|x^j| \le \gamma \eta |x_{\kappa}^1|$ for $j \ne 1$,

 $0 \le t < T$, where γ denotes a constant independent of η , but not always the same. In view of (3.30), this means that (3.28)–(3.29) has a solution $w(t) \ne 0$ such that

$$|w_i^1(t)| \le \gamma \eta s^{\kappa - i} |w_{\kappa}^1(t)|$$
 and $|w^j(t)| \le \gamma \eta s^{\kappa - h_0} |w_{\kappa}^1(t)|$

for $i \neq \kappa, j \neq 1$, respectively. By (3.26), (3.24) has a solution $z(t) \not\equiv 0$ satisfying

$$|z_{i}^{1}(t) - w_{\kappa}^{1}(t)s^{\kappa - i}/(\kappa - i)!| \leq \gamma \eta |w_{\kappa}^{1}(t)|s^{\kappa - i} \quad \text{for } i = 1, \dots, \kappa,$$

$$|z_{i}^{1}(t)| \leq \gamma \eta |w_{\kappa}^{1}(t)|s^{\kappa - i} \quad \text{for } i = \kappa + 1, \dots, h,$$

$$|z^{j}(t)| \leq \gamma \eta |w_{\kappa}^{1}(t)|s^{\kappa - h_{\bullet}} \quad \text{for } j = 2, \dots, g.$$

Finally, by (3.20) and (3.22), (3.19) has a solution $y = (v, v', \ldots, v^{(n-1)})$ satisfying (3.8) if j = 1 and $v = v_{j\kappa}$. This completes the proof of Theorem 3.1.

4. Nonlinear interpolation problems. The results and methods of Lasota and Opial [8], combined with those above, make it possible to treat nonlinear boundary value or interpolation problems associated with certain nonlinear differential equations

$$(4.1) u^{(n)} + a_{n-1}u^{(n-1)} + \cdots + a_0u = f(t, u, u', \ldots, u^{(n-1)}),$$

(4.2)
$$u^{i-1}(t_k) = c_{ik}$$
 for $i = 1, ..., m(k)$ and $k = 1, ..., \nu$,

(4.3)
$$(0=)t_1 < \cdots < t_{\nu} (=T) \text{ and } m(1) + \cdots + m(\nu) = n.$$

ASSUMPTION (B₁). Let a_0, \ldots, a_{n-1} be real numbers satisfying Assumption (A₁) of Theorem 1.1, or the assumption $a_0 = \cdots = a_{n-1} = 0$ of Theorem 2.1, or Assumption (A₂) of Theorem 3.1. Correspondingly, let $0 < T < \infty$ and let $q_0(t), \ldots, q_{n-1}(t)$ be nonnegative, continuous functions on $0 \le t \le T$ such that q(t) defined in (1.4) satisfies (1.5), or q(t) defined in (2.2) satisfies (2.3), or q(t) defined by (3.4) satisfies (3.5).

Assumption (B₂). Let $f(t, y_0, ..., y_{n-1})$ be a continuous function for $0 \le t \le T$ and arbitrary $(y_0, ..., y_{n-1})$ satisfying

$$|f(t, y_0, \ldots, y_{n-1})| \leq \sum_{k=0}^{n-1} q_k(t) |y_k|.$$

THEOREM 4.1. Let a_0, \ldots, a_{n-1} and $q_0(t), \ldots, q_{n-1}(t)$ satisfy Assumption (B₁) and $f(t, y_0, \ldots, y_{n-1})$ satisfy Assumption (B₂). Then the interpolation problem (4.1)–(4.3), in which $c_{11}, \ldots, c_{m(v)v}$ are arbitrary constants, has at least one solution. If, in addition, $f(t, y_0, \ldots, y_{n-1})$ satisfies

$$(4.5) |f(t, y_0, \ldots, y_{n-1}) - f(t, z_0, \ldots, z_{n-1})| \le \sum_{k=0}^{n-1} q_k(t) |y_k - z_k|,$$

then the solution of (4.1)-(4.3) is unique.

Proof. If $\nu = n$ and $m(1) = \cdots = m(n) = 1$, then this theorem follows from [8]. In fact, since T is finite and the main conditions in (1.5), (2.3), (3.5) involve a strict inequality, it is clear that these conditions hold whenever $q_k(t)$ is replaced by a continuous function $Q_k(t)$ satisfying $|Q_k(t)| < q_k(t) + (\eta_* - \eta)/nT$, where $\eta_* = \eta_0$, η_1 , η_2 , respectively. The proof for the general case (4.1)–(4.3) is similar.

APPENDIX

Let $y^k = (y_1^k, \ldots, y_{m(k)}^k)$ be a real Euclidean vector of dimension m(k) and $y = (y^1, \ldots, y^N) = (y_1^1, \ldots, y_{m(N)}^N)$ the corresponding Euclidean vector of dimension $M = m(1) + \cdots + m(N)$. Consider a linear differential equation

$$(1) y' = A(t)y,$$

where A(t) is an $M \times M$ matrix with real continuous entries for $0 \le t < T (\le \infty)$. Let κ be fixed, $1 \le \kappa \le N$. Assume that there exists continuous functions $\psi_1(t) \ge 0$, ..., $\psi_n(t) \ge 0$ and $\alpha_1(t), \ldots, \alpha_{\kappa}(t)$ and $\alpha^{\kappa}(t), \ldots, \alpha^{N}(t)$ such that (1) implies

(2)
$$y^k \cdot y^{k'} \le \alpha_k(t) |y^k|^2 + \psi_k(t) |y^k| \cdot |y|$$
 for $k = 1, ..., \kappa$,

(3)
$$y^k \cdot y^{k'} \ge \alpha^k(t) |y^k|^2 - \psi_k(t) |y^k| \cdot |y|$$
 for $k = \kappa, ..., N$.

Assume that the integrals

(4)
$$\sigma_k(t) = \int_0^t \left[\psi_k(s) + \psi_\kappa(s)\right] \left\{ \exp \int_s^t \left[\alpha_k(\rho) - \alpha^\kappa(\rho)\right] d\rho \right\} ds \text{ for } k = 1, \ldots, \kappa - 1,$$

(5)
$$\sigma_k(t) = \int_t^T \left[\psi_k(s) + \psi_k(s)\right] \left\{ \exp \left[-\int_t^s \left(\alpha^k(\rho) - \alpha_k(\rho)\right] d\rho \right\} ds$$

for
$$k = \kappa + 1, \ldots, N$$
,

exist and are bounded for $0 \le t < T$. Suppose finally that there exist positive constants γ , η with the properties that

(6)
$$\gamma \sigma_k(t) \leq \eta \quad \text{for } 0 \leq t < T, k \neq \kappa,$$

(7)
$$[1+(N-1)\eta^2]^{1/2} \max(1,\eta) < \gamma.$$

THEOREM (*). Under the conditions enumerated above, there exists an $m(1) + \cdots + m(\kappa)$ parameter family of solutions of (1) such that, for $0 \le t < T$, we have $y^{\kappa}(t) \ne 0$ and

(8)
$$|y^{k}(t)| \leq \gamma \sigma_{k}(t)|y^{k}(t)| \quad \text{for } k \neq \kappa.$$

REMARK 1. The proof will show that one obtains a solution of (1) satisfying (8) for which one can assign an arbitrary partial set of initial conditions $y^k(0)$, $k=1,\ldots,\kappa$ such that $y^k(0)\neq 0$ and $|y^k(0)|/|y^k(0)|$ is small for $k=1,\ldots,\kappa-1$.

REMARK 2. Theorem (*) has an analogue if the components of y and the entries of A(t) are complex-valued. It is only necessary to replace $y^k \cdot y^{k'}$ by Re $\bar{y}^k \cdot y^{k'}$ in (2), (3).

REMARK 3. Note that there is no assumption about the signs of the functions α_k , α^k , $\alpha^k - \alpha_k$, $\alpha^k - \alpha_\kappa$.

For particular cases of this theorem, see [2, pp. 284-290; especially, Lemmas 4.1, 4.2 and Exercise 4.3].

Proof. Theorem (*) will be deduced from a result of Ważewski [15]. We shall use the formulation and notation of [2, Theorem 3.1, p 282]. In order to be able to use

this theorem, we shall suppose that $\psi_k(t) > 0$ for 0 < t < T and prove the existence of solutions satisfying $y^k(t) \neq 0$ and (8) for $t_0 \leq t < T$, where t_0 is an arbitrary point of (0, T). The passage to the case $\psi(t) \geq 0$ and $t_0 = 0$ will be clear.

Let Ω be the open (t, y)-set $\Omega = \{(t, y) : 0 < t < T, y \neq 0 \text{ arbitrary}\}$. Define an open subset by

(9)
$$\Omega^0 = \{(t, y) \in \Omega : u_i(t, y) < 0, v_k(t, y) < 0\}$$

for $1 \le k < \kappa < j \le N$, where

$$u_b(t, y) = |y^b|^2 - \gamma^2 \sigma_b^2(t) |y^\kappa|^2 \quad \text{for } b = \kappa + 1, \dots, N,$$

$$v_a(t, y) = |y^a|^2 - \gamma^2 \sigma_a^2(t) |y^\kappa|^2 \quad \text{for } a = 1, \dots, \kappa - 1.$$

Correspondingly, define the subsets of Ω

$$U_b = \{u_b(t, y) = 0, u_j(t, y) \le 0, v_k(t, y) \le 0 \text{ for } 1 \le k < \kappa < j \le N\},\$$

$$V_a = \{v_a(t, y) = 0, u_j(t, y) \le 0, v_k(t, y) \le 0 \text{ for } 1 \le k < \kappa < j \le N\},\$$

where $a=1,\ldots,\kappa-1$ and $b=\kappa+1,\ldots,N$.

If f(t, y) is a function of class C^1 , let $\dot{f}(t, y)$ be the trajectory derivative of f relative to (1); i.e.,

$$f(t, v) = \partial f/\partial t + (\operatorname{grad}_{v} f) \cdot A(t) v.$$

Then we can verify

(10)
$$\dot{v}_a(t, y) > 0$$
 for $(t, y) \in V_a$, $a = 1, ..., \kappa - 1$,

(11)
$$\dot{u}_b(t, y) < 0 \text{ for } (t, y) \in U_b, b = \kappa + 1, \dots, N.$$

For example, in order to obtain (11), note that

$$\dot{u}_b = 2(y^b \cdot y^{b'} - \gamma^2 \sigma_b^2 y^k \cdot y^{k'} - \gamma^2 \sigma_b \sigma_b' |y^k|^2).$$

Using (2) and (3), it is seen that

$$\frac{1}{2}\dot{u}_{b} \geq \alpha^{b}|y^{b}|^{2} - \psi_{b}(t)|y^{b}| \cdot |y| - \gamma^{2}\sigma_{b}^{2}(\alpha_{\kappa}|y^{\kappa}|^{2} + \psi_{\kappa}(t)|y^{\kappa}| \cdot |y|) - \gamma^{2}\sigma_{b}\sigma_{b}'|y^{\kappa}|^{2}.$$

For $(t, y) \in U_b$, we have $|y^b| = \gamma \sigma_b |y^k|$ and

$$|y| \leq |y^{\kappa}| \left(1 + \sum_{j \neq \kappa} \gamma^2 \sigma_j^2\right)^{1/2} \leq c|y^{\kappa}|,$$

where, by (6), $c = [1 + (N-1)\eta^2]^{1/2} > 0$. Thus we have

$$\frac{1}{2}\dot{u}_h \geq \gamma^2 \sigma_h |\gamma^k|^2 \{ (\alpha^b - \alpha_k) \sigma_h - \sigma_h' - (c/\gamma) [\psi_h + \psi_k \gamma \sigma_h] \}.$$

Note that, by (5),

$$\sigma_b' = -(\psi_b + \psi_\kappa) + (\alpha^b - \alpha_\kappa)\sigma_b;$$

hence, for $(t, y) \in U_b$,

$$\frac{1}{2}\dot{u}_b \geq \gamma \sigma_b |y^{\kappa}|^2 \{ (\gamma - c)\psi_b + \psi_{\kappa} (\gamma - c\gamma \sigma_b) \}.$$

Thus $\frac{1}{2}\dot{u}_b > 0$ if $\gamma > c$ and $\gamma > c\gamma\sigma_b$, while $\gamma > c\gamma\sigma_b$ if $\gamma > c\eta$. (Note that $\gamma^* \neq 0$, for

otherwise $(t, y) \in U_b \subset \Omega$ implies that y = 0, which cannot hold for $(t, y) \in \Omega$.) Similarly (10) is verified.

Thus, in the terminology of [2, pp. 281–282], Ω^0 is a (u, v)-subset of Ω and, by a theorem of Ważewski [13] (cf. [2, Lemma 3.1, p. 281]), the set of egress points Ω_e^0 of Ω^0 consists only of strict egress points and

(12)
$$\Omega_{\epsilon}^{0} = \bigcup_{b=\kappa+1}^{N} U_{b} - \bigcup_{a=1}^{\kappa-1} V_{a}.$$

If $0 < t_0 < T$, let $y^1(t_0), \ldots, y^\kappa(t_0)$ be arbitrary vectors satisfying $y^\kappa(t_0) \neq 0$, $|y^a(t_0)| < \gamma \sigma_a(t_0)|y^\kappa(t_0)|$ for $a=1,\ldots,\kappa-1$. Let S be the set of points $(t,y)=(t_0,y^1(t_0),\ldots,y^\kappa(t_0),y^{\kappa+1},\ldots,y^N) \in \Omega$ satisfying $|y^b| \leq \gamma \sigma_b(t_0)|y^\kappa(t_0)|$ for $b=\kappa+1,\ldots,N$. Topologically, S is a ball of dimension $m(\kappa+1)+\cdots+m(N)$ and $S \cap Q_e^0$ is its boundary. Thus $S \cap \Omega_e^0$ is not a retract of S. But $S \cap \Omega_e^0$ is a retract of Ω_e^0 . In fact, a retraction is given by the map $\pi: \Omega_e^0 \to S \cap \Omega_e^0$ defined by $\pi(t,y)=(t_0,z)$, where $z^a=y^a(t_0)$ for $a=1,\ldots,\kappa$ and $z^b=y^b\sigma_b(t_0)|y^\kappa(t_0)|/\sigma_a(t)|y^\kappa|$ for $b=\kappa+1,\ldots,N$. The map π is continuous (since $y\neq 0$, hence $y^\kappa\neq 0$, on Ω_e^0 and $\sigma_b(t)>0$), $\pi(\Omega_e^0)\subset S\cap \Omega_e^0$ and $\pi\mid S\cap \Omega_e^0$ is the identity.

By a theorem of Ważewski [13] (cf. [2, Theorem 3.1, p. 282]), it follows that there exists a point $(t_0, y(t_0)) \in S$ such that the solution of (1) determined by this initial condition is in Ω^0 for $t_0 \le t < T$, i.e., satisfies (8). This proves Theorem (*).

REMARK. It is clear from the proof of this theorem that the conclusion is valid if the functions (4)–(5) and conditions (6)–(7) are replaced by a set of functions $\sigma_a(t)$ and $\sigma_b(t)$, $1 \le a < \kappa < b \le N$, which are nonnegative, continuously differentiable for $0 \le t < T$, and satisfy the following system of differential inequalities

$$\sigma_b' = [\alpha^b(t) - \alpha_k(t)]\sigma_b - \chi_b(t), \qquad \sigma_a' = [\alpha_a(t) - \alpha^k(t)]\sigma_a + \chi_a(t),$$

where χ_a , χ_b are continuous functions satisfying

$$\gamma \chi_k > \left(1 + \gamma^2 \sum_{j \neq \kappa} \sigma_j^2\right)^{1/2} [\psi_k + \psi_\kappa \gamma \sigma_k], \qquad k \neq \kappa.$$

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