

THE DIFFERENTIATION OF PSEUDOINVERSES AND NONLINEAR LEAST  
SQUARES PROBLEMS WHOSE VARIABLES SEPARATE

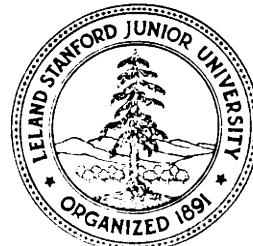
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The Differentiation of Pseudoinverses  
and Nonlinear Least Squares Problems  
Whose Variables-Separate

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## ABSTRACT

For given data  $(t_i, y_i)$ ,  $i=1, \dots, m$ , we consider the least squares fit of nonlinear models of the form

$$F(\underline{a}, \underline{\alpha}; t) = \sum_{j=1}^n g_j(\underline{a}) \varphi_j(\underline{\alpha}; t), \quad \underline{a} \in \mathbb{R}^s, \quad \underline{\alpha} \in \mathbb{R}^k.$$

For this purpose we study the minimization of the nonlinear functional

$$r(\underline{a}, \underline{\alpha}) = \sum_{i=1}^m (y_i - F(\underline{a}, \underline{\alpha}, t_i))^2.$$

It is shown that by defining the matrix  $\{\Phi(\underline{\alpha})\}_{i,j} = \varphi_i(\underline{\alpha}; t_j)$ , and the modified functional  $r_2(\underline{\alpha}) = \| \underline{y} - \Phi(\underline{\alpha}) \Phi^+(\underline{\alpha}) \underline{y} \|_2^2$ , it is possible to optimize first with respect to the parameters  $\underline{\alpha}$ , and then to obtain, a posteriori, the optimal parameters  $\hat{\underline{a}}$ . The matrix  $\Phi^+(\underline{\alpha})$  is the Moore-Penrose generalized inverse of  $\Phi(\underline{\alpha})$ , and we develop formulas for its Fréchet derivative under the hypothesis that  $\Phi(\underline{\alpha})$  is of constant (though not necessarily full) rank. From these formulas we readily obtain the derivatives of the orthogonal projectors associated with  $\Phi(\underline{\alpha})$ , and also that of the functional  $r_2(\underline{\alpha})$ . Detailed algorithms are presented which make extensive use of well-known reliable linear least squares techniques, and numerical results and comparisons are given. These results are generalizations of those of H. D. Scolnik[20].

## 1. Introduction

The least squares fit of **experimental** data is a common tool in many applied sciences and in engineering problems. Linear problems have been well-studied, and stable and efficient methods are available (see for instance: Björck and Golub [3], Golub[8]).

Methods for the nonlinear problems fall mainly in two categories:

(a) general minimization techniques; (b) methods of Gauss-Newton type. The latter type of method takes into consideration the fact that the functional to be minimized is a sum of squares of functions (cf. Daniel [ 5], Osborne [14], Pereyra [15]). The well-known reliable linear techniques have been used mainly in connection with the successive linearization of the nonlinear models. Very recently it has been noticed that by restricting the class of models to be treated, a much more significant use of linear techniques can be made (cf. [2, 9, 12, 13, 17, 20]).

In this paper we consider the following problem. Given data  $(t_i, y_i)$ ,  $i=1, \dots, m$ , find optimal parameters  $\hat{\alpha} = (\hat{\alpha}_1, \dots, \hat{\alpha}_s)^T$ ,  $\hat{\beta} = (\hat{\beta}_1, \dots, \hat{\beta}_k)^T$  that minimize the nonlinear functional

$$(1.1) \quad r(\hat{\alpha}, \hat{\beta}) = \sum_{i=1}^m [y_i - \sum_{j=1}^n g_j(\hat{\alpha}) \varphi_j(\hat{\beta}; t_i)]^2.$$

Throughout this paper a lower case letter in bold face will indicate a column vector, while the same letter with a subscript will indicate a component of the vector. Matrices which are not vectors are denoted by capital letters, and the  $(i,j)$  element of (say) a matrix  $A$  will be indicated by either  $a_{ij}$  or  $\{A\}_{i,j}$ . The transpose of a vector  $\tilde{u}$  is indicated by  $\tilde{u}^T$ . Given a function  $f(t)$ , we shall denote by  $\tilde{f}$  the vector whose components are

$(f(t_1), f(t_2), \dots, f(t_m))^T$ . The scalar product of two vectors  $\underline{u}$  and  $\underline{v}$  is indicated by

$$\langle \underline{u}, \underline{v} \rangle = \underline{v}^T \underline{u} .$$

The only norm which will be used is the Euclidean norm,  $\|\underline{v}\|^2 = \langle \underline{v}, \underline{v} \rangle$ .

Given a matrix  $A$  and a vector  $\underline{b}$ , then we say

$$A\underline{x} \approx \underline{b}$$

if  $\underline{x} = A^+ \underline{b}$  where  $A^+$  is the Moore-Penrose pseudoinverse.

We shall use the symbol  $\mathbf{D}$  for the **Fréchet** derivative of a mapping and  $V$  for the gradient of a functional. We assume that the reader has some familiarity with pseudoinverses and **Fréchet** derivatives and their properties. A useful reference for the pseudoinverse is [19]; for details on the formalism and manipulation of **Fréchet** derivatives, we suggest [6, chapter 8].

Let

$$\{\Phi\}_{i,j} = \Phi_j(\underline{x}; t_i) \quad (i=1, \dots, m; j=1, 2, \dots, n) ,$$

and

$$\underline{g}(\underline{a}) = (g_1(\underline{a}), g_2(\underline{a}), \dots, g_n(\underline{a}))^T .$$

With the given notation, we can rewrite (1.1) as

$$(1.1') \quad r(\underline{a}, \underline{x}) = \|\underline{x} - \Phi(\underline{x}) \underline{g}(\underline{a})\|^2 .$$

Our approach to finding a critical point or a minimum of the functional (1.1') requires two additional hypotheses:

H-1. For any vector  $\underline{b} \in \mathbb{R}^n$ , the system of nonlinear equations

$$(1.2) \quad \underline{g}(\underline{a}) = \underline{b} ,$$

has a solution (not necessarily unique).

H-2. The matrix  $\Phi(\underline{\alpha})$  has constant rank,  $r \leq \min(m, n)$  for  $\underline{\alpha} \in \Omega \subset \mathbb{R}^k$ ,  
 $\Omega$  being an open set containing the desired solution.

Our aim is to be able to deal separately with the parameters  $\underline{\alpha}$ , and then proceed to obtain the parameters  $\underline{a}$ , as it was done in [9, 20] whose results this paper generalize. The reader should also note the independent results obtained by Pérez and Scolnik [17], who in addition deal with nonlinear constraints.

In order to obtain this separation of variables, we consider, as in [9, 17, 20], the modified functional

$$(1.3) \quad r_2(\underline{\alpha}) = \| \underline{y} - \Phi(\underline{\alpha})\Phi^+(\underline{\alpha})\underline{y} \|^2,$$

which will be called the variable projection functional. Once optimal parameters  $\hat{\underline{\alpha}}$  have been obtained by minimizing (1.3), then auxiliary parameters  $\hat{\underline{b}}$  are obtained as  $\hat{\underline{b}} = \Phi^+(\hat{\underline{\alpha}})\underline{y}$ , and finally we take  $\hat{\underline{a}}$  as any solution of the system of equations (1.2).

We shall show in Theorem 2.1 the relationship between critical or minimal points encountered considering the original functional  $r(\underline{a}, \underline{\alpha})$  and those obtained from the functional  $r_2(\underline{\alpha})$  and  $\Phi^+(\underline{\alpha})\underline{y}$ . Both for our proof and for the numerical algorithms of Section 5, we need to develop formulas for the Fréchet derivative of the pseudoinverse of a matrix function. In Section 4, we develop these formulas and obtain the derivatives of the projectors and the Jacobian of the residual vector. The only hypothesis necessary on the rank of the matrix is that it should be constant on an open neighborhood of the point in which the derivative has to be calculated. This is necessary since otherwise the pseudoinverse is not a continuous function, and therefore it could hardly be differentiable. Our proof is coordinate-free. For the full rank case,

similar formulas have been obtained by Fletcher and **Lill** [7] (without proof), by Hanson and Lawson [10], and by **Pérez** and Scolnik [17]. In [7] and [17] this is used to deal with constraints via penalty functions. In [17] the authors choose to work with components, and also obtain a formula for the rank deficient case which is given in terms of the factors of a certain decomposition of the original matrix. Our formulas, besides being coordinate-free and thus much more convenient for algebraic manipulation, are given exclusively in terms of the original matrix, its derivative, and its pseudoinverse. The formula for the rank deficient case seems to be new.

In Section 5 we give a detailed explanation of how to implement the method in an efficient way and in Section 6 we present some numerical examples and comparisons. Extensive use is made of linear least squares techniques.

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2. A class of nonlinear least square-s problems whose parameters separate.

We are going to consider in this paper models of the form:

$$(2.1) \quad n(\underline{a}, \underline{\alpha}; t) = \sum_{j=1}^n g_j(\underline{a}) \varphi_j(\underline{\alpha}; t),$$

where  $\underline{a} \in \mathbb{R}^s$ ,  $\underline{\alpha} \in \mathbb{R}^k$ , and the functions  $g_j$ ,  $\varphi_j$ , are continuously differentiable with respect to  $\underline{a}$ , and  $\underline{\alpha}$  respectively. We shall call the functions  $g_j$  autonomous, to distinguish them from the  $\varphi_j$  which are dependent on  $t$ .

We remark that the parameters  $\underline{a}$  and  $\underline{\alpha}$  form two completely disjoint sets.

The independent variable  $t$  could be a vector itself as in [9, 17].

This requires only small notational changes and we shall not pursue it here.

Given the data  $(t_i, y_i)$ ,  $i=1, \dots, m$ ,  $m \geq s + k$ , our task is to find the values of the parameters  $\underline{a}$ ,  $\underline{\alpha}$ , that minimize the nonlinear functional:

$$(2.2) \quad r(\underline{a}, \underline{\alpha}) = \| \underline{y} - n(\underline{a}, \underline{\alpha}) \|^2 = \sum_{i=1}^m (y_i - n(\underline{a}, \underline{\alpha}; t_i))^2.$$

The approach to the solution of this problem is, as in [9, 17, 20], to modify the functional  $r(\underline{a}, \underline{\alpha})$ , in such a way that consideration of the autonomous parameters  $\underline{a}$  is deferred.

In what follows we shall call  $\Phi(\underline{\alpha})$  the matrix function

$$(2.3) \quad \Phi(\underline{\alpha}) = [\varphi_1(\underline{\alpha}), \dots, \varphi_n(\underline{\alpha})].$$

For each fixed  $\underline{\alpha}$ , the linear operator

$$(2.4) \quad P_{\Phi}(\underline{\alpha}) = \Phi(\underline{\alpha}) \Phi^+(\underline{\alpha}),$$

is the orthogonal projector on the linear space spanned by the columns of the matrix  $\Phi(\underline{\alpha})$ . We shall denote the linear operator  $(I - P_{\Phi}(\underline{\alpha}))$  by  $P_{\Phi}^\perp(\underline{\alpha})$ .

$P_{\Phi}^{\perp}(\underline{\alpha})$  is the projector on the orthogonal complement of the column space of  $\Phi(\underline{\alpha})$ . Similarly,

$$(2.4') \quad \underline{\Phi}^P = \underline{\Phi}^+ \underline{\Phi}$$

is the orthogonal projector on the row space of  $\Phi$ , and  $\underline{\Phi}^{P^{\perp}} = I - \underline{\Phi}^P$ .

When there is no possibility of confusion we shall omit either the matrix subindex or the arguments in projections and functions, or both.

Taking  $\underline{b}$  as a new parameter vector, we consider the following auxiliary model:

$$(2.5) \quad \underline{h}_1(\underline{b}, \underline{\alpha}; t) = \sum_{j=1}^n b_j \varphi_j(\underline{\alpha}; t) .$$

We define similarly the functional  $r_1(\underline{b}, \underline{a}) = \|\underline{y} - \underline{h}_1\|^2$ .

For any given  $\underline{\alpha}$  we have the minimal least squares solution

$$(2.6) \quad \underline{b}^* = \underline{\Phi}^+(\underline{\alpha}) \underline{y} .$$

Thus,

$$(2.7) \quad \min_{\underline{b}} r_1(\underline{b}, \underline{\alpha}) = r_1(\underline{b}^*, \underline{\alpha}) = \|\underline{y} - \underline{\Phi}(\underline{\alpha}) \underline{\Phi}^+(\underline{\alpha}) \underline{y}\|^2 = \|\underline{P}_{\underline{\Phi}(\underline{\alpha})}^{\perp} \underline{y}\|^2 .$$

The modified functional is then the variable projection functional that we mentioned earlier and can be rewritten as:

$$(2.8) \quad r_2(\underline{\alpha}) = \|\underline{P}_{\underline{\Phi}(\underline{\alpha})}^{\perp} \underline{y}\|^2 .$$

Once a critical point (or a minimizer)  $\hat{\underline{\alpha}}$  is found for this functional, then  $\hat{\underline{b}}$  is obtained by replacing  $\underline{\alpha}$  by  $\hat{\underline{\alpha}}$  in (2.6). Finally, by hypothesis H-1,  $\hat{\underline{a}}$  is obtained as any solution of the system of nonlinear equations

$$(2.9) \quad \underline{g}(\underline{a}) = \hat{\underline{b}} .$$

The justification for employing this procedure is given by the following theorem:

Theorem 2.1 Let  $r(\tilde{a}, \tilde{\alpha})$  and  $r_2(\tilde{\alpha})$  be defined as above. We assume that in the open set  $\Omega \subset \mathbb{R}^k$ ,  $\Phi(\tilde{\alpha})$  has constant rank  $r \leq \min(m, n)$ .

(a) If  $\hat{\alpha}$  is a critical point (or a global minimizer in  $\Omega$ ) of  $r_2(\tilde{\alpha})$ , and  $\hat{a}$  satisfies:

$$(2.10) \quad g(\hat{a}) = \Phi^+(\hat{\alpha}) y,$$

then  $\hat{\alpha}$  is a critical point of  $r(\tilde{a}, \tilde{\alpha})$  (or a global. minimizer for  $\tilde{\alpha} \in \Omega$ ) and  $r(\hat{a}, \hat{\alpha}) = r_2(\hat{\alpha})$ .

(b) If  $(\hat{a}, \hat{\alpha})$  is a global minimizer of  $r(\tilde{a}, \tilde{\alpha})$  for  $\tilde{\alpha} \in \Omega$ , then  $\hat{\alpha}$  is a global minimizer of  $r_2(\tilde{\alpha})$  in  $\Omega$  and  $r_2(\hat{\alpha}) = r(\hat{a}, \hat{\alpha})$ .

Furthermore, if there is an unique  $\hat{a}$  among the minimizing pairs of  $r(\tilde{a}, \tilde{\alpha})$ , then  $\hat{a}$  must satisfy (2.10).

We shall postpone the proof of this Theorem until the end of Section 4, where we obtain a convenient expression for the gradient of the functional  $r_2(\tilde{\alpha})$ .

### 3. Algorithmia I. Residual calculation.

One of our main points in the algorithmic part of this paper is to emphasize, when possible and appropriate, the use of stable and efficient linear least squares techniques. Thus it is convenient to review some of the tools and introduce the necessary notation.

If  $Q$  is an orthogonal matrix then, for every vector  $\tilde{z}$ ,

$$\|Q\tilde{z}\| = \|\tilde{z}\|.$$

It is well-known (cf. [8, 10, 18]) that every  $m \times n$  matrix  $\Phi$  ( $m \geq n$ ) of rank  $r \leq n$ , can be orthogonally transformed into "triangular" form, viz., there exist  $Q, Z$  orthogonal, such that

$$(3.1) \quad Q \Phi Z^T = \begin{bmatrix} \tilde{T} & 0 \\ 0 & 0 \end{bmatrix} \equiv T,$$

where  $\tilde{T}$  is an  $r \times r$  upper triangular and nonsingular matrix. Then

$$\Phi^+ = Z^T T^+ Q = Z^T \begin{bmatrix} \tilde{T}^{-1} & 0 \\ 0 & 0 \end{bmatrix} Q,$$

and consequently,

$$P_\Phi = \Phi \Phi^+ = Q^T \begin{bmatrix} I_r & 0 \\ 0 & 0 \end{bmatrix} Q; \quad P_{\Phi^\perp} = Q^T \begin{bmatrix} 0 & 0 \\ 0 & I_{m-r} \end{bmatrix} Q.$$

(Similarly, though we don't use it in our calculations: 

$$\Phi^P = \Phi^+ \Phi = Z^T \begin{bmatrix} I_r & 0 \\ 0 & 0 \end{bmatrix} Z; \quad \Phi^{P^\perp} = Z^T \begin{bmatrix} 0 & 0 \\ 0 & I_{n-r} \end{bmatrix} Z.)$$

Due to the isometric properties of the orthogonal transformation  $Q$ ,

the least squares problem can be expressed as

$$\min_{\tilde{b}} \|\tilde{y} - \tilde{\Phi}\tilde{b}\|^2 = \min_{\tilde{b}} \|\tilde{Q}\tilde{y} - \tilde{\Phi}\tilde{b}\|^2.$$

Calling  $\tilde{y} = \tilde{Q}\tilde{y}$  and partitioning it as  $\tilde{y} = \begin{bmatrix} \tilde{y}_1 \\ \tilde{y}_2 \end{bmatrix}_{(m-r)}$ , we obtain

$$(3.2) \quad \begin{bmatrix} \tilde{b} \\ \tilde{b} \end{bmatrix} = Z^T \begin{bmatrix} \tilde{T}^{-1} \tilde{y}_1 \\ \tilde{0} \end{bmatrix}.$$

A simple computation shows that:

$$(3.3) \quad \|\tilde{y} - \tilde{\Phi}\tilde{b}\|^2 = \|\tilde{P}_{\tilde{\Phi}}^{\perp} \tilde{y}\|^2 = \|\tilde{y}_2\|^2.$$

Therefore, one can evaluate the nonlinear functional  $r_2(\tilde{\alpha})$  of (2.8) for any value of  $\tilde{\alpha}$  in the following way: First the orthogonal matrix  $Q(\tilde{\alpha})$  that is used in the reduction of  $\tilde{\Phi}(\tilde{\alpha})$  is determined; simultaneously,  $\tilde{y} = \tilde{Q}\tilde{y}$  is computed, and finally

$$(3.4) \quad r_2(\tilde{\alpha}) = \|\tilde{y}_2\|^2$$

is evaluated.

For minimization techniques not requiring derivatives this is all that is needed. For iterative techniques using the gradient of the functional or the Jacobian of the residual vector function  $\tilde{P}_{\tilde{\Phi}(\tilde{\alpha})}^{\perp} \tilde{y}$ , we shall provide in the next section formulas which will also be useful in the proof of Theorem 2.1.

#### 4. Fréchet derivatives of pseudoinverses, projectors, and residual vectors.

In this section we develop formulas for the Fréchet derivative of the pseudoinverse of a matrix function. This leads to **expressions** for the derivatives of the associated orthogonal projectors, and for the residual vector function

$$(4.1) \quad \tilde{r}_2(\tilde{\alpha}) = P_{\Phi(\tilde{\alpha})}^{\perp} \tilde{y}.$$

As an aid to those readers not familiar with these concepts, we observe that an  $m \times n$  matrix function  $A(\tilde{\alpha})$  is a nonlinear mapping between the linear space of parameters  $\tilde{\alpha} \in \mathbb{R}^k$  and the space of linear transformations  $\mathcal{L}(\mathbb{R}^n, \mathbb{R}^m)$ . Consequently,  $\mathbf{D}A(\tilde{\alpha})$  will be, for each  $\tilde{\alpha}$ , an element of  $\mathcal{L}(\mathbb{R}^k, \mathcal{L}(\mathbb{R}^n, \mathbb{R}^m))$ . Thus,  $\mathbf{D}A(\tilde{\alpha})$  could be interpreted as a tridimensional tensor, formed with  $k$  ( $m \times n$ ) matrices (slabs), each one containing the partial derivatives of the elements of  $A$  with respect to one of the variables  $\alpha_i$ . Still in another way, each column in the  $k$ -direction is the gradient of the corresponding matrix element.

Since all dimensions involved are different, it will be always clear in the algebraic manipulations how the different vectors, matrices, and tensors interact.

Lemma 4.1. For any  $\tilde{\alpha} \in \Omega$ , an open set of  $\mathbb{R}^k$ , let  $B(\tilde{\alpha})$  be an  $m \times n$  full column rank matrix function, and  $C(\tilde{\alpha})$  an  $n \times m$  full row rank matrix function.

If  $B(\tilde{\alpha})$  and  $C(\tilde{\alpha})$  are Fréchet differentiable in  $\Omega$ , then

$$(4.2) \quad \mathbf{D}(B^+) = -B^+ \mathbf{D}B B^+ + (B^T B)^{-1} \mathbf{D}B^T P_B^{\perp},$$

$$(4.3) \quad \mathbf{D}(C^+) = -C^+ \mathbf{D}C C^+ + C^T P^{\perp} \mathbf{D}C^T (C C^T)^{-1}.$$

proof. Since  $B$  has full column rank, then  $B^+ = (B^T B)^{-1} B^T$ , and

$$D(B^+) = D(B^T B)^{-1} B^T + (B^T B)^{-1} D B^T .$$

But,

$$D(B^T B)^{-1} = -(B^T B)^{-1} D(B^T B) (B^T B)^{-1} .$$

Therefore,

$$(4.4) \quad D(B^+) = (B^T B)^{-1} [ D B^T - D(B^T B) B^+ ] .$$

Developing  $D(B^T B)$  and regrouping, we obtain (4.2). Since  $C^T$  has full column rank, (4.3) follows readily from (4.2). ■

Since  $P_A(\underline{\alpha}) = AA^+$ ,  $P_A^\perp(\underline{\alpha}) = I - AA^+$ , it follows that

$$(4.5) \quad D P_A = D_A A^+ + {}_A D(A^+) ,$$

and

$$(4.6) \quad D P_A^\perp = -D P_A .$$

If  $A(\underline{\alpha})$  has full column rank, then from (4.5) and Lemma 4.1 we obtain

$$(4.7) \quad D P_A P_A^\perp D_A A^+ + (P_A^\perp D_A A^+)^T .$$

Similarly, if  $A$  has full row rank:

$$(4.8) \quad D_A P = A^+ D_A A^P + (A^+ D_A A^P)^T .$$

We shall prove now that formulas (4.7) and (4.8) are valid in the rank deficient case. For this purpose we shall prove first an auxiliary Lemma, and then obtain the derivative of the pseudoinverse of an arbitrary matrix function.

Let  $A(\underline{\alpha})$  be an  $m \times n$  matrix function, Fréchet differentiable,

and with constant rank  $r \leq \min(m, n)$ , on an open set  $\Omega \subset \mathbb{R}^k$ . Let  $B(\alpha)$  be a maximal set of independent columns of  $A(\alpha)$  in  $\Omega$ , and let  $C = \overset{\sim}{B+A}$ . It is well-known (see, for instance, [16]): (1)  $C$  has full row rank, (2)  $A = BC$ , (3)  $A^+ = C+B^+$ . Due to our hypothesis,  $B(\alpha)$  can be formed with the same columns of  $A(\alpha)$  on a neighborhood of every  $\alpha \in \Omega$ , Other useful relations that we shall use below are

$$AA^+ = P_A = BB^+ = P_B; B^+ P_A^\perp = 0;$$

$$CA^+ = B^+ AA^+ = B^+; P_A^\perp B = 0; P_A^\perp A^+ = 0.$$

Lemma 4.2. With  $A$ ,  $B$ , and  $C$  defined as above, the following formula is valid in  $\Omega$ :

$$(4.9) \quad B D B^+ P_A^\perp = (D_A A^+, P_A^\perp).$$

Proof: From Lemma 4.1 we get

$$\begin{aligned} B D B^+ &= -P_B D B B^+ + (D B B^+)^T P_B^\perp \\ &= -P_A D B B^+ + (D B B^+)^T P_A^\perp. \end{aligned}$$

Therefore,

$$(4.10) \quad B D B^+ P_A^\perp = (D B B^+)^T P_A^\perp,$$

On the other hand, since  $A = BC$ ,

$$D_A A^+ = D_B C A^+ + B D C A^+ = D_B B^+ + B D C A^+.$$

Thus,

$$P_A^\perp D_A A^+ = P_A^\perp D_B B^+,$$

or

$$(D_A A^+)^T P_A^\perp = (D_B B^+)^T P_A^\perp,$$

and this last expression together with (4.10) proves the Lemma. ■

Theorem 4.3.  $\subseteq \mathbb{R}^k$  be an open set and for  $\alpha \in \Omega$  let  $A(\alpha)$  be an  $m \times n$  Fréchet differentiable matrix function having fixed rank  $r \leq \min(m, n)$ .

Then for any  $\alpha \in \Omega$ :

$$(4.11) \quad D A^+ = -A^+ D_A A^+ + A^+ A^{+T} D_A^T P_A^\perp + P_A^\perp D_A^T A^+ A^{+T}.$$

proof: With B and C as above, we have that

$$D A^+ = D(C^+ B^+) = D C^+ B^+ + C^+ D_B^+,$$

and hence by (4.3)

$$D A^+ = -C^+ D_C C^+ B^+ + C^+ P_A^\perp D_A^T C^+ B^+ + C^+ D_B^+,$$

since

$$(C C^T)^{-1} = C^T C^+.$$

Substituting

$$D_C = D(B^+ A) = D_B^+ A + B^+ D_A, \quad C^+ B^+ = A^+, \quad C^+ P_A^\perp = P_A^\perp,$$

in the last expression we get:

$$(4.12) \quad \mathbf{D} \mathbf{A}^+ = -\mathbf{A}^+ \mathbf{D} \mathbf{A} \mathbf{A}^+ + \mathbf{C}^+ \mathbf{D} \mathbf{B}^+ - \mathbf{C}^+ \mathbf{D} \mathbf{B}^+ \mathbf{A} \mathbf{A}^+ \mathbf{A}^P \mathbf{D} \mathbf{C}^T \mathbf{C}^+ \mathbf{A}^+$$

$$= -\mathbf{A}^+ \mathbf{D} \mathbf{A} \mathbf{A}^+ + \mathbf{C}^+ \mathbf{D} \mathbf{B}^+ \mathbf{P}_A^+ \mathbf{A}^P \mathbf{D} \mathbf{C}^T \mathbf{C}^+ \mathbf{A}^+$$

But,

$$(4.13) \quad \mathbf{A}^P \mathbf{D} \mathbf{C}^T \mathbf{C}^+ \mathbf{A}^+ = \mathbf{A}^P \mathbf{D} \mathbf{A}^T \mathbf{A}^+ \mathbf{A}^+ + \mathbf{A}^P \mathbf{A}^T \mathbf{D} \mathbf{B}^+ \mathbf{C}^+ \mathbf{A}^+$$

$$= \mathbf{A}^P \mathbf{D} \mathbf{A}^T \mathbf{A}^+ \mathbf{A}^+ ,$$

since  $\mathbf{A}^P \mathbf{A}^T = 0$ .

Substituting (4.13) into (4.12) and using the relationship  $\mathbf{C}^+ \mathbf{D} \mathbf{B}^+ \mathbf{P}_A^+ = \mathbf{A}^+ \mathbf{B} \mathbf{D} \mathbf{B}^+ \mathbf{P}_A^+ = \mathbf{A}^+ \mathbf{A}^+ \mathbf{D} \mathbf{A}^T \mathbf{P}_A^+$ , given by Lemma 4.2, we finally obtain the desired result. ■

Corollary 4.3. Let  $\mathbf{A}(\alpha)$  be as in Theorem 4.2. Then, for any  $\alpha \in \Omega$ .

$$(4.14a) \quad \mathbf{D} \mathbf{P}_A = \mathbf{D}(\mathbf{A} \mathbf{A}^+) = \mathbf{P}_A^T \mathbf{D} \mathbf{A} \mathbf{A}^+ + (\mathbf{P}_A^T \mathbf{D} \mathbf{A} \mathbf{A}^+) ,$$

$$(4.14b) \quad \mathbf{D} \mathbf{A}^P = \mathbf{D}(\mathbf{A}^+ \mathbf{A}) = \mathbf{A}^+ \mathbf{D} \mathbf{A} \mathbf{A}^P + (\mathbf{A}^+ \mathbf{D} \mathbf{A} \mathbf{A}^P)^T .$$

Proof: Obvious. I

From this result it is now easy to derive an expression for the gradient of the functional  $r_2(\alpha)$  (see (2.8)), provided the matrix  $\Phi(\alpha)$  has constant rank on an open neighborhood of the point in which the gradient is calculated.

In fact:

$$(4.15) \quad r_2(\alpha) = \| \mathbf{P}_{\Phi(\alpha)}^T \mathbf{y} \|^2 = \langle \mathbf{P}_{\Phi(\alpha)}^T \mathbf{y}, \mathbf{P}_{\Phi(\alpha)}^T \mathbf{y} \rangle ,$$

and

$$\frac{1}{2} \nabla r_2(\underline{\alpha}) = -\underline{y}^T \underline{P}^{\frac{1}{2}} \underline{D}^{\Phi} \underline{\Phi}^+ \underline{\Phi}^T \underline{D}^{\Phi} \underline{P}^{\frac{1}{2}} \underline{y}.$$

Since  $\underline{P}^{\frac{1}{2}} \underline{\Phi}^T = 0$ , we finally obtain:

$$(4.16) \quad \frac{1}{2} \nabla r_2(\underline{\alpha}) = -\underline{y}^T \underline{P}^{\frac{1}{2}} \underline{D}^{\Phi} \underline{\Phi}^+ \underline{y}.$$

Now we have the elements for proving Theorem 2.1.

Proof of Theorem 2.1.

From (2.2) we have that  $r(\underline{a}, \underline{\alpha}) = \|\underline{y} - \underline{\Phi}(\underline{\alpha})\underline{g}(\underline{a})\|^2$ .

Therefore,

$$(4.17) \quad \frac{1}{2} \nabla r(\underline{a}, \underline{\alpha}) = -(\underline{y} - \underline{\Phi}\underline{g})^T (\underline{D}\underline{\Phi} \underline{\Phi}^+ \underline{y} + \underline{\Phi}\underline{D}\underline{g}).$$

Assume now that  $\hat{\underline{\alpha}}$  is a critical point of  $r_2(\underline{\alpha})$ , and that  $\hat{\underline{a}}$  satisfies

$$(4.18) \quad \underline{g}(\hat{\underline{\alpha}}) = \underline{\Phi}^+(\hat{\underline{\alpha}})\underline{y}.$$

Then,

$$(4.19) \quad \begin{aligned} \frac{1}{2} \nabla r(\hat{\underline{a}}, \hat{\underline{\alpha}}) &= -(\underline{P}^{\frac{1}{2}} \underline{y})^T (\underline{D}\underline{\Phi} \underline{\Phi}^+ \underline{y} + \underline{\Phi}\underline{D}\underline{g}) \\ &= \frac{1}{2} \nabla r_2(\hat{\underline{\alpha}}), \end{aligned}$$

since  $\underline{y}^T \underline{P}^{\frac{1}{2}} \underline{\Phi} \underline{D}\underline{g} = 0$ . Thus  $(\hat{\underline{a}}, \hat{\underline{\alpha}})$  is a critical point of  $r(\underline{a}, \underline{\alpha})$ ,

Assume now that  $\hat{\underline{\alpha}}$  is a global minimizer of  $r_2(\underline{\alpha})$  in  $\Omega$ , and  $\hat{\underline{a}}$  satisfies (4.18). Then clearly,  $r(\hat{\underline{a}}, \hat{\underline{\alpha}}) = r_2(\hat{\underline{\alpha}})$ . Assume that there exists  $(\underline{a}^*, \underline{\alpha}^*)$ ,  $\underline{\alpha}^* \in \Omega$ , such that  $r(\underline{a}^*, \underline{\alpha}^*) < r(\hat{\underline{a}}, \hat{\underline{\alpha}})$ . Since for any  $\underline{\alpha}$  we have  $r_2(\underline{\alpha}) \leq r(\underline{a}, \underline{\alpha})$ , then it follows that  $r_2(\underline{\alpha}^*) \leq r(\underline{a}^*, \underline{\alpha}^*) < r(\hat{\underline{a}}, \hat{\underline{\alpha}}) = r_2(\hat{\underline{\alpha}})$ , which is a contradiction to the fact that  $\underline{\alpha}$  was a global

minimizer of  $r_2(\underline{a})$  in  $\Omega$ . Therefore  $(\hat{\underline{a}}, \hat{\underline{\alpha}})$  is a global **minimizer** of  $r(\underline{a}, \underline{\alpha})$  in  $\Omega$  and part (a) of the Theorem is **proved**.

Conversely, suppose that  $(\hat{\underline{a}}, \hat{\underline{\alpha}})$  is a global minimizer of  $r(\underline{a}, \underline{\alpha})$  in  $\Omega$ , then as above

$$r_2(\hat{\underline{\alpha}}) \leq r(\hat{\underline{a}}, \hat{\underline{\alpha}}).$$

Now let  $\underline{a}^*$  be a solution of  $g(\underline{a}) = \Phi^+(\hat{\underline{\alpha}})y$ .

Then we have

$$r_2(\hat{\underline{\alpha}}) = r(\underline{a}^*, \hat{\underline{\alpha}}) \leq r(\hat{\underline{a}}, \hat{\underline{\alpha}}),$$

but since  $(\hat{\underline{a}}, \hat{\underline{\alpha}})$  was a global minimizer we must have equality. If there was an unique  $\underline{a}$  among the minimizers of  $r(\underline{a}, \underline{\alpha})$ , then  $\underline{a}^* \equiv \hat{\underline{a}}$ . We still have to show that  $\hat{\underline{\alpha}}$  is a global minimizer of  $r_2(\underline{\alpha})$ . Assume that it is not. Thus, there will be  $\bar{\underline{\alpha}} \in \Omega$ , such that  $r_2(\bar{\underline{\alpha}}) < r_2(\hat{\underline{\alpha}})$ . Let  $\bar{\underline{a}}$  be a solution of  $g(\underline{a}) = \Phi^+(\bar{\underline{\alpha}})y$ . Then  $r_2(\bar{\underline{\alpha}}) = r(\bar{\underline{a}}, \bar{\underline{\alpha}}) < r_2(\hat{\underline{\alpha}}) = r(\hat{\underline{a}}, \hat{\underline{\alpha}})$ , which is a contradiction to the fact that  $(\hat{\underline{a}}, \hat{\underline{\alpha}})$  was a global minimizer of  $r(\underline{a}, \underline{\alpha})$ . I

5. Algorithmia II. Detailed implementation of the Gauss-Newton-Marquardt algorithm.

We shall now explain in detail how to apply the results of Section 4 to the Marquardt modification of the Gauss-Newton iterative procedure; we make extensive use of linear least squares techniques. We shall include an economical implementation of the Marquardt algorithm devised earlier by Golub (see also [11, 14]).

We define the vector

$$\underline{r}_2(\underline{\alpha}) = P_{\Phi}^{\frac{1}{2}}(\underline{\alpha}) \underline{y} .$$

The generalized Gauss-Newton iteration-with step control for the nonlinear least squares problem

$$(5.1) \quad \min_{\underline{\alpha}} \underline{r}_2(\underline{\alpha}) = \min_{\underline{\alpha}} \left\| \underline{r}_2(\underline{\alpha}) \right\| = \min_{\underline{\alpha}} \left\| P_{\Phi}^{\frac{1}{2}}(\underline{\alpha}) \underline{y} \right\|$$

is given by

G.N. Starting from an arbitrary  $\underline{\alpha}^0$  :

$$(5.2) \quad \underline{\alpha}^{\ell+1} = \underline{\alpha}^{\ell} - t_{\ell} [\mathbf{D}\underline{r}_2(\underline{\alpha}^{\ell})]^+ \underline{r}_2(\underline{\alpha}^{\ell}) , \quad (\ell=0, \dots) .$$

The parameters  $t_{\ell} > 0$ , which control the size of the step, are used to prevent divergence. Usually  $t_a = 1$ , unless  $\underline{r}_2(\underline{\alpha}^{\ell+1}) > \underline{r}_2(\underline{\alpha}^{\ell})$ , in which case  $t_{\ell}$  is reduced. Another use of the parameters  $t_{\ell}$  is to minimize  $\underline{r}_2(\underline{\alpha}^{\ell+1})$  along the direction  $[\mathbf{D}\underline{r}_2(\underline{\alpha}^{\ell})]^+ \underline{r}_2(\underline{\alpha}^{\ell})$ .

Marquardt's modification calls for the introduction of a sequence of non negative auxiliary parameters  $v_{\ell} > 0$ .

G.N.M. Define

$$K_{\ell} \equiv \begin{bmatrix} \mathbf{D}\underline{r}_2(\underline{\alpha}^{\ell}) \\ \hline v_{\ell} F_{\ell} \end{bmatrix} , \quad \underline{z}_{\ell} \equiv \begin{bmatrix} \underline{r}_2(\underline{\alpha}^{\ell}) \\ \hline 0 \end{bmatrix} \}_{n} ,$$

where for each  $\ell$ ,  $F_{\ell}$  is the upper triangular Cholesky factor of an  $n \times n$

symmetric positive definite matrix  $M_\ell$ . Then the Gauss-Newton-Marquardt iteration is given by

$$\tilde{\alpha}^{\ell+1} = \tilde{\alpha}^\ell - K_{\ell+1}^+ \tilde{r}_\ell, \quad \ell \geq 0.$$

Reasons for this modification are well-known. For more details and an interesting study of the convergence of this method we refer to [14]. We wish to make explicit now the "two-stage" orthogonal factorization" given in [11] and [14], in order to show how to take advantage of the special structure of the problem.

$$\text{Calling } \tilde{h} = \tilde{\alpha}^{\ell+1} - \tilde{\alpha}^\ell, \quad DP = D\tilde{r}_\ell(\tilde{\alpha}^\ell) = DP_{\Phi}^\ell \tilde{y}$$

and dropping the superscript  $\ell$  from here on in, one step of the Marquardt algorithm is equivalent to solving the linear least squares problem

$$Kh \approx \begin{bmatrix} -\tilde{r}_2(\tilde{\alpha}) \\ \tilde{\alpha} \end{bmatrix}.$$

In the first stage of the orthogonal factorization of  $K$  an  $m \times n$  orthogonal matrix  $Q$  is chosen so that

$$Q DP = R_1 = \begin{bmatrix} \text{Hatched} \\ 0 \\ \text{O} \end{bmatrix}^n = \begin{bmatrix} R'_1 \\ \text{O} \end{bmatrix}.$$

Thus,

$$\begin{aligned} \begin{bmatrix} Q & | & 0 \\ \hline 0 & | & I_n \end{bmatrix} \cdot \begin{bmatrix} DP \\ vF \end{bmatrix} &\equiv Q_1 \begin{bmatrix} DP \\ vF \end{bmatrix} = \begin{bmatrix} R'_1 \\ vF \end{bmatrix}, \\ Q_1 \begin{bmatrix} -\tilde{r}_2(\tilde{\alpha}) \\ \tilde{\alpha} \end{bmatrix} &\equiv \begin{bmatrix} \tilde{r} \\ \tilde{\alpha} \end{bmatrix}. \end{aligned}$$

$R'_1$  and  $\tilde{r}$  are saved for future use.

In the second stage we choose an  $(m+n) \times (m+n)$  orthogonal matrix  $Q_2$  to reduce

$$A \equiv \begin{bmatrix} R'_1 \\ \hline \sqrt{F} \end{bmatrix} \text{ to "triangular" form.}$$

For this purpose we shall use successive Householder transformations as in [3], from where we adopt the notation.\*

On reducing the first column of  $A$ , which is of the form:

$$\tilde{a}_1^{(1)} = \begin{bmatrix} a_{11}^{(1)} \\ 0 \\ \vdots \\ 0 \\ \hline \mu \\ 0 \\ \vdots \\ 0 \end{bmatrix} \quad m, \quad \mu = \sqrt{f_{11}}$$

we use  $Q^{(1)} = I - \beta_1 \tilde{u}_1^{(1)} \tilde{u}_1^{(1)T}$ ,

where

$$u_1^{(1)} = \text{sign}(a_{11}^{(1)}) (\sigma_1 + |a_{11}^{(1)}|) ,$$

$$\sigma_1 = (a_{11}^2 + \mu^2)^{\frac{1}{2}} ,$$

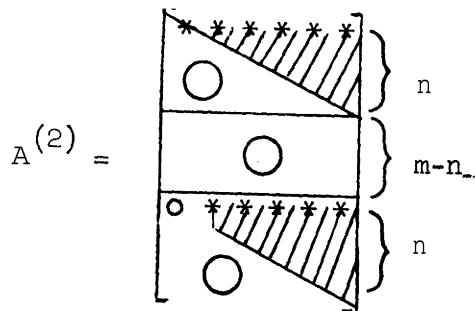
$$u_{m+1}^{(1)} = \mu ,$$

$$u_1^{(1)} = 0, \text{ otherwise,}$$

$$\beta_1 = 2/\tilde{u}_1^{(1)T} \tilde{u}_1^{(1)} .$$

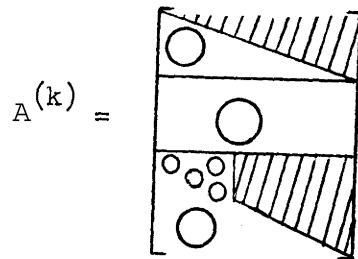
Now we observe that when  $Q^{(1)}$  is applied to a vector, any component corresponding to a zero component of  $\tilde{u}_1^{(1)}$  is left unchanged. In particular,

the band of zeros in  $A$  is preserved. Thus, in this first step we only need to transform the elements of rows number 1 and  $m+1$ . Consequently,  $A^{(2)} = Q^{(1)} A$  will have the schematic form:



where the asterisks indicate the modified elements.

It is now clear that at step  $k$ ,  $A^{(k)}$  will have the form



The matrix  $A^{(k+1)}$ ,  $k = 1, \dots, n$ , is obtained as follows:

$$i) \sigma_k = ((a_{kk}^{(k)})^2 + \sum_{i=1}^k (a_{m+i,k}^{(k)})^2)^{\frac{1}{2}},$$

$$ii) \beta_k = (\sigma_k (\sigma_k + |a_{kk}^{(k)}|))^{-1},$$

$$iii) u_i^{(k)} = 0 \text{ for } i < k, k+1 \leq i \leq m, m+k < i;$$

$$u_k^{(k)} = \text{sign}(a_{kk}^{(k)}) (\sigma_k + |a_{kk}^{(k)}|);$$

$$u_i^{(k)} = a_i^{(k)}, \quad m+1 \leq i \leq m+k.$$

$$iv) \quad y^T = \beta_k u^{(k)T} A^{(k)} ,$$

$$y_j = \beta_k [u_k^{(k)} a_{kj}^{(k)} + \sum_{i=1}^k a_{m+i,k}^{(k)} a_{m+i,j}^{(k)}] , \quad j=k+1, \dots, n .$$

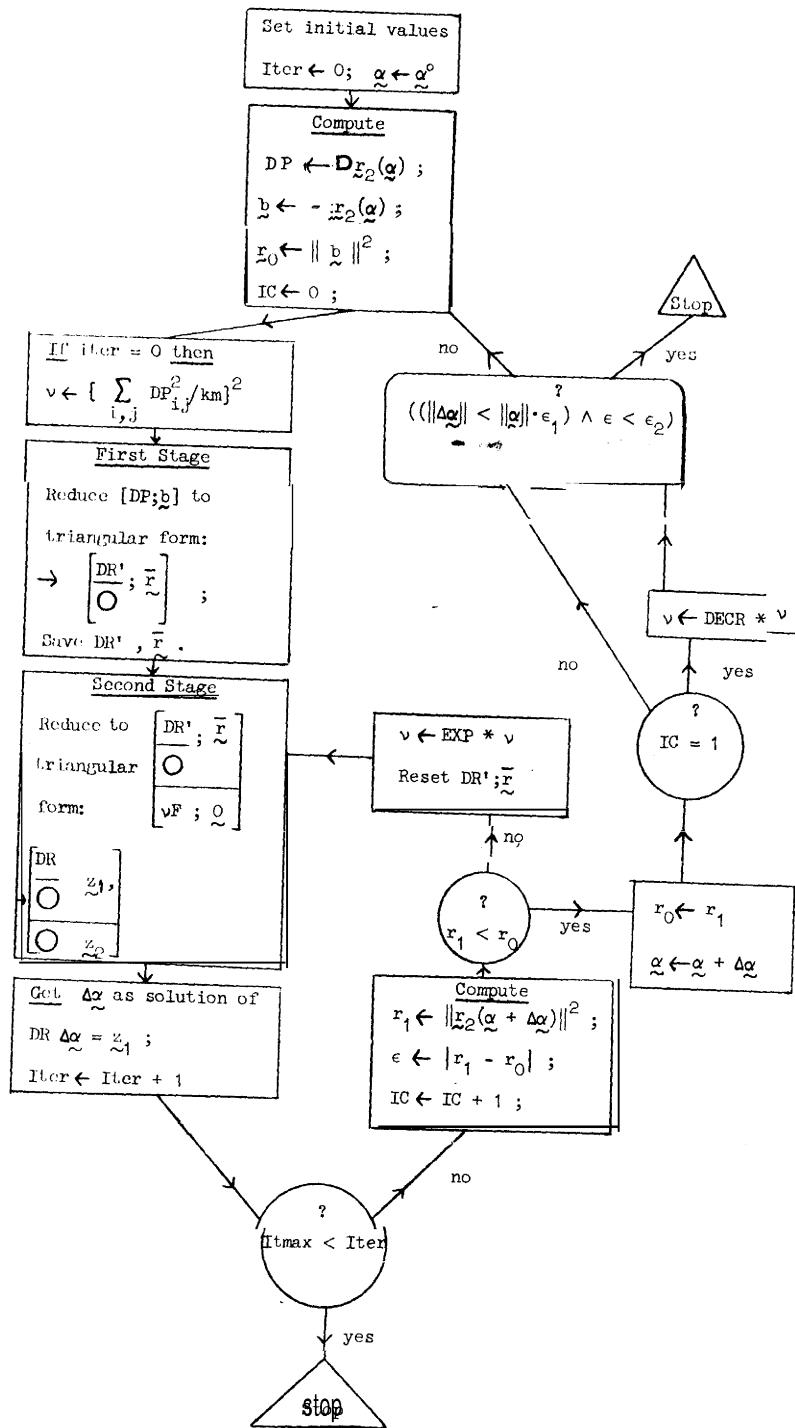
Finally,

$$v) \quad a_{ij}^{(k+1)} = a_{ij}^{(k)} - u_1^{(k)} y_j , \quad \begin{matrix} i=k; m+1, \dots, m+k; \\ j=k+1, \dots, n; \end{matrix}$$
$$a_{kk}^{(k+1)} = \text{sign} (a_{kk}^{(k)}) \sigma_k . \quad - ,$$

These formulas are similar to those given in [3], but are modified to take advantage of the structure of the matrix A .

Osborne's version of Marquardt's algorithm, modified for our present problem, is presented in the detailed flow-chart of Figure 1. The parameters DECR and EXP are the factors by which  $\nu$  is either decreased or increased.

Figure Marquardt's algorithm (Osborne [14])



We will evaluate  $\mathbf{D}\underline{x}_2(\underline{\alpha}) = \mathbf{D}\underline{P}_\Phi^\perp(\underline{\alpha})\underline{y}$  for a given  $\underline{\alpha}$ , according to

$$(5.4) \quad \mathbf{D}\underline{P}_\Phi^\perp(\underline{\alpha})\underline{y} = -(\underline{P}_\Phi^\perp \mathbf{D} \mathbf{a}) \Phi^+ \underline{y} - \Phi^+ (\underline{P}_\Phi^\perp \mathbf{D}\Phi)^\top \underline{y},$$

which is readily obtained from (4.14a).

In many applications, each component function  $\varphi_j$  depends only upon a few of the parameters  $\{\alpha_t\}_{t=1}^k$ , and therefore its derivatives with respect to the other parameters will vanish. Those vanishing derivatives will produce  $m$ -columns of zeros in the tensor  $\mathbf{D}\Phi$ . In order to avoid a waste of storage and useless computation with zeros it is convenient to introduce from the outset the  $k \times n$  incidence matrix  $E = (e_{jt})$ . This matrix will be defined as follows:

$e_{jt} = 1$  iff parameter  $\alpha_t$  appears in function  $\varphi_j$ ;

$e_{jt} = 0$  otherwise.

We shall also call  $p$  the number of nonzero derivatives in  $\mathbf{D}\Phi$ :  $p = \sum_{t,j} e_{jt}$ .

The nonzero derivative vectors can then be stored sequentially in a bidimensional array  $B(m \times p)$ . In our implementation we chose to store the nonzero  $m$ -columns varying first the index corresponding to the different differentiations, and then that corresponding to the different functions. This information can then be decoded for use in algebraic manipulations by means of the incidence matrix  $E$ ,

We now introduce some notation in order to describe the compressed storage of the nonzero columns of the tensor  $\mathbf{D}\Phi$  in a more explicit fashion. We define, for  $t = 1, \dots, k$ ,

$S_t = \{ \text{set of ordered indices for which } e_{jt} \neq 0, j=1, \dots, n \};$

$$\psi_{tj}(\underline{\alpha}) = \frac{\partial \varphi_j(\underline{\alpha})}{\partial \alpha_t} \quad , \quad j=1, \dots, n \quad , \quad t=1, \dots, k \quad .$$

We write the matrix B in partitioned form

$$B = [B_1, B_2, \dots, B_k] \quad ,$$

where

$$B_t = [\psi_{tj_1}, \psi_{tj_2}, \dots, \psi_{tj_t}]_{j_i \in S_t} \quad .$$

A step-by-step description of the computation of  $D\Phi^{-1}(\underline{\alpha})$  follows.

We assume that the rank of  $\Phi(\underline{\alpha})$  is computationally determined and equal to  $r \leq \min(m, n)$ .

- a) Compute  $\Phi(\underline{\alpha})$ ,  $D\Phi(\underline{\alpha})$ .
- b) Form the  $m \times (n+p+1)$  array

$$G \equiv [\Phi(\underline{\alpha}) ; \underline{y} ; D\Phi(\underline{\alpha})] = [A ; \underline{y} ; B] \quad .$$

- c) Obtain the complete orthogonal factorization of A (cf. Section 3):

$$QAZ^T = T = \begin{bmatrix} \widetilde{T} & 0 \\ 0 & 0 \end{bmatrix} ; \quad \widetilde{T} = \begin{bmatrix} \text{triangular} \\ \text{circles} \end{bmatrix}_{r \times r} \quad .$$

$$\text{Also } \underline{v} = Q\underline{y} ; \quad C = QB$$

( $\widetilde{T}$ ,  $\underline{v}$ , and C will be stored in the array G). Note again that (see Section 3):

$$\begin{bmatrix} \widetilde{T} \\ \underline{v} \end{bmatrix} = \begin{bmatrix} 0 & 0 \\ \hline 0 & I_{n-m+1} \end{bmatrix} Q \quad .$$

- d) Get the intermediary values:

$$\bar{D} = \begin{bmatrix} 0 & 0 \\ \hline 0 & I_{m-r} \end{bmatrix} \cdot C \quad (\text{i.e., Remember that the nonzero information of } \bar{D} \text{ is stored in the last } p \text{ columns and last } m-r \text{ rows of } G);$$

$$\tilde{x} = A^+ \tilde{y} = Z^T \begin{bmatrix} \tilde{T}^{-1} v_1 \\ \hline \tilde{v}_2 \end{bmatrix} \quad (v = \begin{bmatrix} v_1 \\ \hline v_2 \end{bmatrix} \}^r_{m-r}) .$$

$$e) \quad \underset{n \times k}{U} = (P^T D \Phi)^T \tilde{y} = D \Phi^T Q^T \begin{bmatrix} 0 & 0 \\ \hline 0 & I_{m-r} \end{bmatrix} Q \tilde{y} = \bar{D}^T \tilde{y} = \bar{D} \otimes \begin{bmatrix} 0 \\ \hline v_2 \end{bmatrix}$$

(transposition in the tensor  $D \Phi$  refers to transposition within the "slabs" corresponding to the different derivatives, and must be interpreted adequately when decoding the information from the compressed storage array  $G$ ; the appropriate ALGOL-60 code for computing  $U$  with our storage convention would be (assuming that  $C = QB$  is stored in the same place  $B$  is :

```

n1 ← n+1 ;
L ← n1 ;
for t ← 1 step 1 until k do
  for j ← 1 step 1 until n do
    if E[j,t] = 0 then U[j,t] ← 0 else
      begin L ← L+1 ; acum ← 0 ;
      for i ← n1 step 1 until m do
        acum ← acum + G[i,L] X G[i,n1] ;
      U[j,t] ← acum
      end ; )
  
```

f) Compute  $S = Z \cdot U$  .  
 $\underset{n \times k}{}$

Solve the  $k \times r$  lower triangular systems:

$\tilde{T}^T W = \tilde{S}$  , where  $\tilde{S}$  contains the first  $r$  rows of  $S$  .  
 $\underset{r \times k}{}$

Store  $W$  in the first  $r$  rows of the  $m \times k$  array  $B$ . Compute  $\bar{D} \otimes x$  and store the nonzero information in the last  $m-r$  rows of  $B$ .

g) Finally, the  $m \times k$  matrix  $B$  is obtained as:

$$B \leftarrow D P_{\Phi(\alpha)}^{-1} y = -Q^T \begin{Bmatrix} W \\ \bar{D} \otimes x \end{Bmatrix} = -Q^T B .$$

We emphasize the systematic use made of the triangular orthogonal decomposition of the matrix  $\Phi(\alpha)$ . We also warn the reader about the correct interpretation of the algebraic operations in which any tridimensional tensor intervene, as we exemplified in (e).

## 6. Numerical experiments.

We have implemented three different algorithms based on the developments of the previous sections for the case  $\mathbf{g}(\mathbf{a}) = \mathbf{a}$  and rank  $\Phi = n$ .

The methods minimize the variable projection functional  $r_2(\mathbf{g}) = \|P_{\Phi}^{\perp}(\mathbf{g})\mathbf{g}\|^2$  first, in order to obtain the optimal parameters  $\hat{\mathbf{g}}$ , and then complete the optimization according to our explanation in Section 2. The algorithms differ in the procedure used for the minimization of  $r_2(\mathbf{g})$ .

A1. Minimization without derivatives. We use PRAXIS, a FORTRAN version of a 'program developed by R. Brent [4], who very kindly made it available to us. All that PRAXIS essentially requires from the user is the value of the functional for any  $\mathbf{g}$ . This is computed using the results of Section 3. In fact, the user has only to give code for filling the matrix  $\Phi$  for any  $\mathbf{g}$ , and our program will effect the triangular reduction and so on. It turns out that many times (see the examples) the models have some terms which are exclusively linear, i.e., functions  $\varphi_j$  which are independent of  $\mathbf{g}$ . Those functions produce columns in  $\Phi(\mathbf{g})$  which are constant throughout the process. If they are considered first, then it is possible to reduce them once and for all, saving the repetition of computation. This is done in our program.

### A2. Minimization by Gauss-Newton with control of step (see (5.2)).

The user is required to provide the incidence matrix  $E$  and the array of functions  $\varphi_j$  and non-vanishing partial derivatives:  $G$ . See Section 5 for a more detailed description.

### A3. Minimization by Marquardt's modification, as explained in Section 5 with

$F_{\alpha} \equiv I$ . User supplied information is the same as in A2.

Test problems. Problems 1 and 2 are taken from Osborne [14], where the necessary data can be found.

P1. Exponential fitting. The model is of the form:

$$n_1(a, \alpha; t) = a_1 + a_2 e^{-\alpha_1 t} + a_3 e^{-\alpha_2 t}.$$

The functions  $\varphi_j$  are obviously  $\varphi_1(\alpha; t) \equiv 1$ ,  $\varphi_{j+1}(\alpha; t) = e^{-\alpha_j t}$ ,  $j=1,2$ .

So the different constants, in the notation of Section 2 are:  $n=3$ ,  $s=3$ ,  $k=2$ .

For the problem considered,  $m=33$ . The number of constant functions: NCF = 1. The number of non-vanishing partial derivatives:  $p=2$ .

In Table I we compare our results for methods A1, A2, A3, and those obtained by minimizing the full functional  $r(a, \alpha)$ .

P2. Fitting Gaussians with an exponential background.

$$n_2(a, \alpha; t) = a_1 e^{-\alpha_1 t} + a_2 e^{-\alpha_2(t-\alpha_5)^2} + a_3 e^{-\alpha_3(t-\alpha_6)^2} + a_4 e^{-\alpha_4(t-\alpha_7)^2}.$$

The functions  $\varphi_j$  are:

$$\varphi_1(\alpha; t) = e^{-\alpha_1 t}; \varphi_j(\alpha; t) = e^{-\alpha_j(t-\alpha_{j+3})^2}, \quad j=2,3,4.$$

Thus:  $n=4$ ,  $s=4$ ,  $k=7$ ,  $m=65$ ,  $p=7$ .

Results for this problem appear in Table II.

P3. Iron Mössbauer Spectrum with two sites of different electric field gradient and one single line [21].

The model here is the following:

$$n_3(a, \alpha; t) = a_1 + a_2 t + a_3 t^2$$

$$\begin{aligned}
 & - a_1 \left[ \frac{1}{1 + \left| \frac{\alpha_1 + 0.5\alpha_2 - t}{\alpha_3} \right|^2} + \frac{1}{1 + \left| \frac{\alpha_1 - 0.5\alpha_2 - t}{-\alpha_3} \right|^2} \right] \\
 & - a_5 \left[ \frac{1}{1 + \left| \frac{\alpha_4 + 0.5\alpha_5 - t}{\alpha_6} \right|^2} + \frac{1}{1 + \left| \frac{\alpha_4 - 0.5\alpha_5 - t}{\alpha_6} \right|^2} \right] \\
 & - a_6 \left[ \frac{1}{1 + \left| \frac{\alpha_7 - t}{\alpha_8} \right|^2} \right]
 \end{aligned}$$

Clearly,  $\phi_j(\alpha; t) = t^j$ ,  $j=1,2,3$ ; and  $\phi_4$ ,  $\phi_5$ ,  $\phi_6$  are the functions inside the square brackets.

Here:  $n=6$ ,  $k=8$ ,  $NCF=3$ ,  $p=8$ ,  $m=188$ ,  $s=6$ .

For this example we wish to thank Dr. J. C. Travis of NBS who kindly supplied the problem and results from his own computer program.

Comparisons are offered in Table III.

The qualitative behavior of the three different minimization procedures used in our computation follows the pattern that have been expounded in recent comparisons (Bard [1]). Gauss-Newton is fastest whenever it converges from a good initial estimate. As is shown in the fitting of Gaussians (Table II), if the problem is troublesome, then a more elaborate strategy is called for. Brent's program has the advantage of not needing derivatives, which in this case leads to a big simplification. On the other hand, it is a very conservative program which really tries to obtain rigorous results. This, of course, can lead to a long search in cases where it is not entirely justified.

As a consequence of our Theorem 2.1, and of our numerical experience, we strongly recommend, even in the case when our procedure is not used, to obtain

initial values for the linear parameters when  $g_j(\tilde{a}) = a_j$  by setting

$\tilde{a}^0 = \Phi^+(\tilde{a}^0)y$ . This is done in our program for the full functional and in the program of Travis with excellent results.

The computer times shown in Table I and Table II correspond to the CPU times (execution of the object code) on an IBM 360/50. All calculations were performed in long precision; viz. 14 hexadecimal digits in the mantissa of each number. We compare the results of minimizing the reduced functional when the Variable Projection (VP) technique is used with that of minimizing the full functional (FF) for various minimization algorithms. In order to eliminate the coding aspect, we have used essentially the same code for minimizing the two functionals. The only difference was in the subroutine DPA which computes in both cases the Jacobian of the residual vector.

In the FF approach, the subroutine DPA computed the  $m \times (n+k)$  matrix B as follows: the first n columns consisted of the vectors  $\Phi_j(\tilde{a})$  while the remaining columns were the partial derivatives

$$\frac{\partial}{\partial a_\ell} (y - \Phi(\tilde{a})\tilde{a}) = - \sum_{j=1}^n a_j \frac{\partial \Phi_j(\tilde{a})}{\partial a_\ell}, \quad (\ell=1,2,\dots,k).$$

These derivatives were constructed using the same information provided by the user subroutine ADA. We also obtained from DPA in the FF case, the automatic initialization of the linear parameters, viz.  $\tilde{a}^0 = \Phi^+(\tilde{a}^0)y$ .

For the numerical examples given here, the cost per iteration was somewhat higher for the VP functional. However, we see that in some cases there has been a dramatic decrease in the number of iterations; this has been observed previously (cf. [12]). Thus, in these cases the total computing time is much more favorable for the VP approach. This was especially true for all three

methods of minimization when the exponential fit was made and when Marquardt's method was used in the Mössbauer spectrum problem,

For the Mossbauer spectrum problem, we used two sets of initial values. We used those given by Travis [21], (say)  $\underline{\beta}^0$ , and also  $\hat{\underline{\beta}}^0 \approx \underline{\beta}^0 \pm 0.05 \underline{\beta}^0$ . For  $\underline{\beta}^0$ , the value of the functional is  $3.04467 \times 10^8$  while for  $\hat{\underline{\beta}}^0$ , the value of the functional is  $6.405 \times 10^8$ ; the final estimates of the parameters yielded a residual sum of squares less than  $3.0444 \times 10^8$ . When Brent's method was used on the full functional, the method did not seem to converge, but for the reduced functional, Brent's method converged reasonably well. In fact, after twenty minutes Brent's algorithm applied to the full functional with  $\underline{\beta}^0$  did not achieve the desired reduction in the functional.

The results we have obtained in minimizing the full functional for the first two problems using the Marquardt method, and those of problem 3 with Newton's method and  $\underline{\beta}^0$ , are consistent with the results reported by Osborne and Travis.

From a rough count of the number of arithmetic operations (function and derivative evaluation per step are the same for both procedures, so that the work they do can be disregarded), it seems that for almost no combination of the parameters (m, n, k, p) the VP procedure will require fewer operations per iteration than the FF procedure. It is an open problem then to determine a priori under what conditions the VP procedure will converge more quickly than the FF procedure when minimization algorithms using derivatives are used.

Another important problem is that of stability. The numerical stability of the process and of the attained solution must be studied. By insisting on the use of stable linear techniques, we have tried to achieve an overall numerically stable procedure for this nonlinear situation. Since the standards

of stability for non-linear problems are ill-defined at' this time, it is hard to say whether we have succeeded in obtaining our goal.

Table I

Exponential fit.

Method	Functional	Number of Function Evaluations	Number of Derivative Evaluations	Time (seconds)
$A_1$	FF	1832	—	191.00
	VP	100	—	9.00
$A_2$	FF	11	11	5.05
	VP	4	4	3.20
$A_3$	FF	32	26	12.55
	VP	4	4	3.12

$$r(\hat{a}, \hat{\alpha}), r_2(\hat{\alpha}) \leq 0.5465 \times 10^{-4}$$

$\sim \sim$

Table II

Gaussian fit.

Method	Functional	Number of Function Evaluations	Number of Derivative Evaluations	Time (seconds)
$A_3$	FF	11	9	23.35
	VP	10	8	26.82

$$r(\hat{a}, \hat{\alpha}), r_2(\hat{\alpha}) \leq 0.048$$

Methods  $A_1$  and  $A_2$  were either slowly convergent or non-convergent.

Table III

Mössbauer Iron Spectrum.

Method	Functional	Initial Values	Number of Function Evaluations	Number of Derivative Evaluations	Time (seconds)
$A_1$	FF	$\beta^0$	65	0	*
	VP	$\beta^0$			70.00
$A_2$	FF	$\beta^0$	4	4	34.34
	VP	$\beta^0$	4	4	41.64
	FF	$\beta^0$	7	7	52.27
	VP	$\beta^0$	6	6	59.60
$A_3$	FF	$\beta^0$	16	16	118.22
	VP	$\beta^0$	3	3	35.35
	FF	$\beta^0$	18	18	130.50
	VP	$\beta^0$	6	6	61.92

$$r(\hat{a}, \hat{\alpha}), r_2(\hat{\alpha}) \leq 3.0444 \times 10^8$$

$$(\hat{\beta}^0 = (80, 49, 5, 81, 24, 9.5, 100, 4)^T)$$

\* Did not converge in finite amount of time.

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Key Words

Pseudoinverse

Nonlinear Least Squares

**Fréchet** Derivative

Projectors

Orthogonalization

FORTRAN Program

```

C*****DRIVER PROGRAM FOR USE OF VARPRO IN PROBLEM 1: FITTING OF TWO
C      EXPONENTIALS AND ONE CONSTANT TERM.
C      IMPLICIT REAL*8(A-H,O-Z)
C      DIMENSION Y(200),T(200),ALF(20),AC(20)
C      EXTERNAL ADA
C      CALL LEER(N,M,KG,NCFUN,Y,T,ALF)
C      CALL VARPRO(N,M,KG,NCFUN,Y,T,ALF,AC,ADA)
C      CALL EXIT
C      END

```

C---

```

SUBROUTINE VARPRO(N,M,KG,NCFUN,Y,T,ALF,AC,ADA)
IMPLICIT REAL*8(A-H,O-Z)
COMMON A(200,20),AA(200,10),E(20,20),B(220,20),UKK(200),
* BETA(20),P
INTEGER P
DIMENSION UK1(20),BET1(20),Z(20),DR(20,20),ZPR(200),DEL(20)
* ,ALF(KG),ALF1(20),AC(20),Y(M),T(M)
EXTERNAL ADA

```

NONLINEAR LEAST SQUARES PROGRAM FOR LINEAR COMBINATIONS OF NONLINEAR FUNCTIONS.

WRITTEN IN FORTRAN 4 - LEVEL G. IN THIS SUBROUTINE THERE ARE WRITE STATEMENTS USING UNIT 3 AS OUTPUT. THAT UNIT NUMBER IS INSTALLATION DEPENDENT.

MINIMIZATION BY OSBORNE-MARQUARDT ALGORITHM (OR GAUSS-NEWTON WITH STEP CONTROL BY MAKING THE SMALL CHANGES INDICATED IN THE SECOND LINE AFTER INSTRUCTION Labeled 5, AND AFTER LABEL 61).

SEE 'THE DIFFERENTIATION OF PSEUDOINVERSES AND NONLINEAR LEAST SQUARES PROBLEMS WHOSE VARIABLES SEPARATE' BY GENE H. GOLUB AND V. PEREYRA, STANFORD U. TECHN. REP. 261, MARCH 1972.

M = NUMBER OF OBSERVATIONS.

N = NUMBER OF FUNCTIONS.

KG = NUMBER OF NONLINEAR VARIABLES.

NCFUN = NUMBER OF CONSTANT FUNCTIONS, I.E. FUNCTIONS PHI WHICH DO NOT DEPEND UPON ANY PARAMETERS ALPHA. THEY SHOULD APPEAR FIRST.

Y = M-VECTOR OF OBSERVATIONS.

T = M-VECTOR OF INDEPENDENT VARIABLE,

(= (N\*KG) INCIDENCE MATRIX. E(I,J) = 1 IFF VARIABLE J APPEARS IN FUNCTION I. P = SUM OF E(I,J).

ALF = KG-VECTOR OF INITIAL VALUES. ON OUTPUT IT WILL CONTAIN THE OPTIMAL VALUES OF THE NONLINEAR PARAMETERS.

AC = N-VECTOR OF LINEAR PARAMETERS (OUTPUT).

\*\*\*\*\*

CONTINUE

THE USER MUST PROVIDE A SUBROUTINE THAT FOR GIVEN ALF WILL EVALUATE THE FUNCTIONS PHI AND THE I-TH PARTIAL DERIVATIVES D PHI(I)/D ALF(J), AT THE SAMPLE POINTS T. THE VECTOR SAMPLED FUNCTION PHI(I) SHOULD BE STORED IN THE I-TH COLUMN OF THE (M X (P+N+1)) MATRIX A. THE NONZERO DERIVATIVES COLUMN VECTORS SHOULD BE STORED SEQUENTIALLY IN THE MATRIX A STARTING IN THE COLUMN N+2. IF ITER=0 (THE FIRST TIME THIS SUBROUTINE IS CALLED), THE MATRIX SHOULD BE FILLED. WITH THIS MATRIX THE STORAGE OF THE DERIVATIVES IS EXPLAINED IN THE FOLLOWING CODE:

L = N+1

DO 10 J=1,KG

```

C      DO 10 I=1,N
C      IF (E(I,J))10,10,11
11  L=L+1
C      DO 10 K=1,M
C      A(K,L) *= D PHI(I) / D ALF(J) (T(K))
10  CONTINUE
C      THE N+1-TH COLUMN OF A IS RESERVED FOR THE VECTOR OF OBSERVATIONS Y.
C      THE SUBROUTINE HEADING SHOULD BE: ( ISEL = 0 : FUNCT. AND DER. MUST B E COMP
C      ISEL = -1 : ONLY FUNCTIONS MUST BE COMP. ISEL = 1: ONLY DER. NECESSARY)

C      SUBROUTINE ADA(N,M,KG,A,E,ITER,P,T,ALF,ISEL)
C
C      (ITER IS AN ITERATION COUNTER PROVIDED BY VARPRO).
C      IF IS ASSUMED THAT THE MATRIX PHI (ALPHA) HAS ACWAYS FULL COLUMN RANK
C ****
ITER=0
C*****THE THREE FOLLOWING PARAMETERS ARE USED IN THE CONVERGENCE TEST (BETWEEN
C INSTRUCTIONS NUMBER 200 AND 400): EPS1 IS A RELATIVE TOLERANCE FOR
C DIFFERENCE BETWEEN TWO CONSECUTIVE RESIDUALS; ITMAX IS THE MAXIMUM
C THE SIZE OF THE CORRECTION. EPS2 IS A RELATIVE TOLERANCE FOR THE
C NUMBER OF FUNCTION AND DERIVATIVE EVALUATIONS ALLOWED.
ITMAX=50
EPS1=10D-4
EPS2=50D-6
C*****
KG1=KG+1
DO 10 I=1,M
10  A(I,N+1)=Y(I)
2  CALL DPA(N,M,KG,NCFUN,ITER,ITER,R,Y,T,ALF,ADA)
CT=1.0D0
IF (ITER.EQ.0)
*WRITE(3,104)ITER,R
IC=0
IF (ITER)3,5,3
5  CONTINUE
XNU=0.
C*****
C***** IF GAUSS-NEWTON IS DESIRED REMOVE THE NEXT FOUR (4) STATEMENTS (SEE
C ALSO LABELED 61) .
DO 4 I=1,M
DO 4 J=1,KG
4  XNU=XNU+B(I,J)**2
XNU= DSQRT(XNU/(M*KG))
WRITE(3,105)XNU
C*****REDUCTION OF B TO TRI ANGULAR FORM,
3  DO 30 I=1,KG
SGMA=0
DO 11 I 1=I ,M
11  SGMA=SGMA+B(I,I)**2
SGMA=DSQRT(SGMA)
IF (B(I,I))12,12,13
12  SG=-1.
GO TO 14

```

```

13  SG=1.
14  UK1(I)=SG*(SGMA+DABS(B(I,I)))
    BET1(I)=1./(SGMA*(SGMA+DABS(B(I,I))))
    B(I,I)=-SG*SGMA
    I1=I+1
23  DO 15 I2=I1,KG1
    ACUM=UK1(I)*B(I,I2)
    DO 16 I3=I1,M
16  ACUM=ACUM+B(I3,I)*B(I3,I2)
15  Z(I2)=BET1(I)*ACUM
    DO 30 J=I1,KG1
    B(I,J)=B(I,J)-UK1(I)*Z(J)
    DO 30 J2=I1,M
30  B(I2,J)=B(I2,J)-B(I2,I)*Z(J)      ----
C****SAVE TRIANGULAR FORM AND Z*.
    DO 40 I=1,KG
    DO 40 J=I,KG
40  DR(I,J)=B(I,J)
    DO 41 I=1,M
41  ZPK(I)=B(I,KG1)
C****REDUCTION. SECOND PHASE.
50  IF(XNU .EQ. 0.0)GO TO 300
    DO 60 K=1,KG
    K1=K+1
    B(M+K,K)=XNU
    DO 42 J=K1,KG1
42  B(M+K,J)=0.0
    SGMA=B(K,K)**2
    DO 51 J=1,K
51  SGMA=SGMA+B(J+M,K)**2
    SGMA=DSQRT(SGMA)
    IF(B(K,K)>52,52,53
52  SG=-1.
    GO TO 54
53  SG=1.
54  UK1(K)=SG*(SGMA+DABS(B(K,K)))
    BET1(K)=1./(SGMA*(SGMA+DABS(B(K,K))))
    B(K,K)=-SG*SGMA
    DO 55 J=K1,KG1
    ACUM=UK1(K)*B(K,J)
    IF(K .EQ. 1)GO TO 55
    K2=K-1
    DO 56 I=1,K2
    IM=I+M
    ACUM=ACUM+B(IM,K)*B(IM,J)
55  Z(J)=BET1(K)*ACUM
    DO 57 J=K1,KG1
    B(K,J)=B(K,J)-UK1(K)*Z(J)
    DO 57 I=1,K
    IM=I+I
57  B(MI,J)=B(MI,J)-B(MI,K)*Z(J)
80  QUITINUE
C****SOLVE FOR DELTA-ALF.
300  NZ=KG-1
    DEL(KG)=B(KG,KG1)/B(KG,KG)

```

```

ALF1(KG)=ALF(KG)+DEL(KG)
DO 58 I=1,N2
  I1=KG-I
  I2=I1+1
  ACUM=B(I1,KG1)
  DO 59 J=I2,KG
    59 ACUM=ACUM-B(I1,J)*DEL(J)
    DLL(I1)=ACUM/B(I1,I1)
    58 ALF1(I1)=ALF(I1)+DEL(I1)
C*****GET NEW RESIDUAL.
  310 ITER=ITER+1
  ISEL=-1
  WRITE(3,103)ITER
  WRITE(3,800)(ALF1(I),I=1,KG)
  IF(ITER .GT. ITMAX) GO TO 400
  DO 900 I=1,M
  900 A(I,N+1)=Y(I)
  CALL DPA(N,M,KG,NCFUN,ITER,ISEL,R1,Y,T,ALF1,ADA)
  IC=IC+1
  WRITE(3,107)IC,R1
  IF(R-R1)61,60,60
  61 CONTINUE
C*****IF GAUSS-NEWTON IS DESIRED REMOVE THE C FROM THE NEXT SIX (6) STATEMENTS
C  IF(XNU)110,111,110
C111  TT=0.5*TT
C  IF(TT .LT. 5.0-4)GO TO 400
C  DO 112 I=1,KG
C112  ALF1(I)=ALF(I)+TT*DEL(I)
C  GO TO 310
C*****
  110 XNU=1.5*XNU
  WRITE(3,106)XNU
C*****RETRIEVE TRIANGULAR FORM OF FIRST PHASE.
  DO 62 I=1,KG
  DO 62 J=I,KG
  62 B(I,J)=DR(I,J)
  DO 63 I=1,M
  63 B(I,KG1)=ZPR(I)
  GO TO 50
  60 EPS=R-R1
  R=R1

  ACC=0.
  DAC=0.
  DO 65 I=1,KG
  ALF(I)=ALF1(I)
  ACC=ACC+ALF(I)**2
  65 DAC=DAC+DEL(I)**2
C*****IF IC IS GREATER THAN 1 THEN NU HAS BEEN INCREASED DURING THIS
C  ITERATION.
  IF(IC .EQ. 1) XNU=0.5*XNU
  WRITE(3,200)IC,XNU
  ACC=DSQRT(ACC)
  DAC=DSQRT(DAC)
  AC1=DAC/ACC
  WRITE(3,108)AC1

```

```

11 ((DAB .LE. ACC*EPS1 .AND. EPS .LE. R*EPS2))
* GO TO 400
120 ISEL=1
130 GO TO 2
400 N1=N-1
AC(N)=A(N,N+1)/A(N,N)
IF (N .EQ. 1) GO TO 135
DO 130 I=1,N1
I1=N-I
I2=I1+1
ACUM=A(I1,N+1)
DO 120 J=I2,N
120 ACUM=ACUM-A(I1,J)*AC(J)
130 AC(I1)=ACUM/A(I1,I1)
135 WRITE(3,209)
      WRITE(3,210)(AC(I),I=1,N)
      WRITE(3,215)(ALF(I),I=1,KG)
      WRITE(3,209)
500 RETURN
103 FORMAT(1H0,'    ITER=',I3,'      PARAMETERS')
104 FORMAT(1H0,'  RESIDUAL',I5,D15.7)
105 FORMAT(1H0,'  NU=',D15.7)
106 FORMAT(1H0,'  N U WAS INCREASED TO',D15.7)
137 FORMAT(1H0,I5,'  NEW RESIDUAL',D15.7)
103 FORMAT(1H0,'  THE NORM OF THE RELATIVE CORRECTION IS=',D15.3)
200 FORMAT(1H0,1S,'  NU IS',D15.7)
209 FORMAT(1H0,50('*'))
210 FORMAT(1H0,'  WEIGHTS//(4D15.7)')
215 FORMAT(1H0,'  NONLINEAR PARAMETERS//(4D15.7)')
800 FORMAT(1H0,4020.10)
END

```

```

C
C-----+
C
C      SUBROUTINE DPA(N,M,KG,NCFUN,ITER,ISEL,R,Y,T,ALF,ADA)
C*****COMPUTATION OF THE DERIVATIVE OF THE VARIABLE PROJECTION.
C
C      IMPLICIT REAL*8(A-H,O-Z)
C      COMMON A(200,20),AA(200,10),E(20,20),B(220,20),UKK(200),
C      * BETA(20),P
C      INTEGER P
C      DIMENSION ALF(KG),Z(120),X(20),U(20,20),Y(M),T(M)
C      EXTERNAL ADA
C      CALL ADA(N,M,KG,A,E,ITER,P,T,ALF,ISEL)
C      N1=N+1
C      N2=1
C      IF(ISEL.GT.0)GO TO 111
C      IF(ITER.GT.0)N2=NCFUN+1
C      DO 110 I=1,M
C      DO 110 J=N2,N
C      110 AA(I,J)=A(I,J)
C*****REDUCTION OF A TO TRIANGULAR FORM, COMPUTATION OF V=QY, AND
C      SELECTIVE COMPUTATION OF QB ACCORDING TO VALUE OF ISEL.
111 110 DO 30 I=1,N
      I11=I+1
      IF(ISEL.GT.0) GO TO 22

```

```

IF (ITER .GT. 0 .AND. I .LE. NCFUN) GO TO 0 7
SGMA=0.
DO 1 1 I1=I,M
11 SGMA=SGMA+A(I1,I)**2
SGMA=DSQRT(SGMA)
IF(A(I,I))12,12,13
12 SG=-1.
GO TO 14
7 I1=NCFUN+1
13 GO TO 20
14 SG=1.
UKK(I)=SG*(SGMA+DABS(A(I,I)))
BETA(I)=1. / (SGMA*(SGMA+DABS(A(I,I))))
A(I,I)=-SG*SGMA
I1=I11
8 IF (ISEL)20,21,22
20 NN=N1
GO TO 23
21 NN=N1+P
GO TO 2 3
22 NN=N1+P
I1=N+2
23 DO 1 5 I2=I1,NN
ACUM=UKK(I)*A(I,I2)
DO 1 4 I3=I11,M
16 ACUM=ACUM+A(I3,I)*A(I3,I2)
15 Z(I2)=BETA(I)*ACUM
DO 1 7 J=I1,NN
A(I,J)=A(I,J)-UKK(I)*Z(J)
DO 1 7 I2=I11,M
17 A(I2,J)=A(I2,J)-A(I2,I)*Z(J)
30 CONTINUE
IF (ISEL .GT. 0) GO TO 50
R=0.
DO 40 I=N1,M
40 R=R+A(I,N1)**2
IF (ISEL .LT. 0) RETURN
*****D-SNAKE IS CONTAINED NOM IN A(I,J), I=N+1, . . . , M; J=N+2, . . . , N+P+1.
C COMPUTATION OF X.
50 N2=N-1
X(N)=A(N,N1)/A(N,N)
IF (N .EQ. 1) GO TO 310
DO 300 I=1,N2
I1=N-I
I2=I1+1
ACUM=A(I1,N1)
DO 200 J=I2,N
200 ACUM=ACUM-A(I1,J)*X(J)
300 X(I1)=ACUM/A(I1,I1)
*****COMPUTATION OF U.
310 L=N1
DO 60 J=1,KG
DO 60 I=1,N
IF(E(I,J))70,70,71
70 U(I,J)=0.

```

```

      GO TO 60
71  L=L+1
      ACUM=0.
      DO 600 K=N1,M
600  ACUM=ACUM+A(K,L)*A(K,N1)
      U(I,J)=ACUM
      GO CONTINUE
C****COMPUTATION OF W (STORED IN UPPER PART OF B).
      DO 80 J=1,KG
      B(I,J)=U(I,J)/A(1,1)
      DO 80 I=2,N
      ACUM=U(I,J)
      I1=I-1
      DO 7 9 L=1,II
7 9  ACUM=ACUM-A(L,I)*B(L,J)
      B(I,J)=ACUM/A(I,I)
C****COMPUTATION OF D-SNAKE* X (STORED IN LOWER PART OF B).
      DO 9 0 I=N1,M
      L=N1
      DO 9 0 J=1,KG
      ACUM=0.
      0 0 900 K=1,N
      IF(E(K,J))900,900,92
      3 2 L=L+1
      ACUM=ACUM+A(I,L)*X(K)
900  CONTINUE
      90 B(I,J)=ACUM
C****FINALLY, DPA(ALF)*Y IS PRODUCED AS QT B.
      DO 9 5 K1=1,N
      K=N-K1+1
      DO 93 I=1,KG
      K2=K+1
      ACUM=UKK(K)*B(K,I)
      DO 94 J=K2,M
3 4  ACUM=ACUM+A(J,K)*B(J,I)
      93 Z(I)=BETA(K)*ACUM
      DO 96 J=1,KG
      B(K,J)=B(K,J)-UKK(K)*Z(J)
      DO 9 6 I=K2,M
9 6  B(I,J)=B(I,J)-A(I,K)*Z(J)
      9 5 CONTINUE
C****COMPUTATION OF ETA=ORTOGONAL COMPONENT OF Y, RESPECT OF A.
      DO 1 2 0 I=1,M
      ACUM= Y(I)
      DO 119 J=1,N
119 ACUM=ACUM-AA( I,J)*X(J)
      120 B( I,KG+1)=ACUM
      RETURN
      END
C
C-----
C
      SUBROUTINE ADA(N,M,KG,A,E,ITER,P,T,ALF,ISEL)
C      OSBORNE'S EXPONENTIAL FITTING.TWO EXPONENTIALS AND CONSTANT TERM.
      IMPLICIT REAL*8(A-H,O-Z)

```

```

      INTEGER P
      DIMENSION A(200,20),E(20,20),ALF(KG),T(M)
      L=0
      IF(ITER .GT. 0)GO TO 5
C*****IN THIS CASE THE INCIDENCE MATRIX E IS:
C
C      
$$\begin{Bmatrix} 0 & 0 \\ 1 & 0 \\ 0 & 1 \end{Bmatrix}$$

C
C      E(1,1)=0.
C      E(1,2)=0.
C      E(2,1)=1.
C      E(2,2)=0.
C      E(3,1)=0.
C      E(3,2)=1.
C
      P=2
      DO 4 I=1,M
      4 A(I,1)=1.0D0
      5 IF(ISEL .GT. 0)GO TO 16
      DO 1 0 I=1,M
      A(I,2)=DEXP(-ALF(L+1)*T(I))
      1 0 A(I,3)=DEXP(-ALF(L+2)*T(I))
      IF(ISEL)14,15,16
      1 b DO 17 I=1,M
      A(I,5)=-T(I)*DEXP(-ALF(L+1)*T(I))
      1 7 A(I,6)=-T(I)*DEXP(-ALF(L+2)*T(I))
      14 RETURN
      15 DO 20 I=1,M
      A(I,5)=-T(I)*A(I,2)
      2 0 A(I,6)=-T(I)*A(I,3)
      RETURN
      END

```

```

SUBROUTINE LEER(N,M,KG,NCFUN,Y,T,ALF)
IMPLICIT REAL*8(A-H,O-Z)
DIMENSION Y(200),T(200),ALF(20)
C*****LEER READS THE DATA.SEE FORMATS 100,102.
1  READ(1,100,END=500)N,M,KG,NCFUN,(T(I),Y(I),I=1,M)
100 FORMAT(4I5/(2D15.7))
      WRITE(3,101)N,M,KG,NCFUN,(I,T(I),Y(I),I=1,M)
101 FORMAT(1H1,' NON LINEAR LEAST SQUARES PROBLEM'//' NUMBER Of - FUNC
      *TIONS=' ,13,3X , 'NUMBER OF OBS
      *ERVATIONS=' ,I3// ' NUMBER OF VARIABLES =' ,I3,'NUMBER OF CONSTANT
      * FUNCTIONS=' ,I3// ' '
      *T(I)           Y(I)'//(I5,2D20.7))
      N1=1
      READ(1,102)(ALF(I),I=N1,KG)
102 FORMAT(4D20.7)
      WRITE(3,103)(ALF(I),I=N1,KG)
103 FORMAT(1H0,' INITIAL NONLINEAR PARAMETERS'//(4D20.7))
      WRITE(3,104)
104 FORMAT( 1H0,50('*'))

```

500 RETURN  
CALL EXIT  
END

## NON LINEAR LEAST SQUARES PROBLEM

NUMBER OF FUNCTIONS= 3 NUMBER OF OBSERVATIONS= 33

NUMBER OF VARIABLES = 2 NUMBER OF CONSTANT FUNCTIONS= 1

1	T(I)	Y(I)
1	0.0	0.84400000 0 0
2	0.1000000D 02	0.90800000 0 0
3	0.2000000D 02	0.93200000 00
4	0.3000000D 02	0.93600000 00
5	0.4000000D 02	0.92500000 00
6	0.5000000D 02	0.90800000 00
7	0.60000000 02	0.88100005 00
8	0.7000000D 0 2	0.85000000 0 0
9	0.8000000D 02	0.81800000 00
10	0.90000000 02	0.7840000D 0 0
11	0.1 000000D 0 3	0.75100000 00
12	0.11000000 0 3	0.7180000D 0 0
13	0.12000000 03	0.6850000D 00
14	0.13000000 03	0.6580000D 00
15	0.1400000D 0 3	0.6280000D 0 0
16	0.13000000 0 3	0.6030000D 0 0
17	0.16000000 0 3	0.58000000 00
18	0.17000000 03	0.5580000D 00
19	0.18000000 03	0.5380000D 00
20	0.1900000D 0 3	0.5220000D 00
21	0.20000000 0 3	0.50600000 00
22	0.21000000 0 3	0.4900000D 00
23	0.22000000 0 3	0.4780000D 00
24	0.23000000 0 3	0.4670000D 00
25	0.2400000D 0 3	0.4570000D 00
26	0.2500000D 0 3	0.44800000 00
27	0.26000000 0 3	0.4380000D 00
28	0.2700000D 0 3	0.4310000D 00
29	0.2800000D 0 3	0.4240000D 00
30	0.2900000D 0 3	0.4200000D 00
31	0.300~0000 03	0.41400000 00
32	0.31000000 0 3	0.4110000D 00
33	0.3200000D 0 3	0.4060000D 00

## INITIAL NONLINEAR PARAMETERS

0.1000000D-01 0.2000000D-01

\*\*\*\*\*

RESIDUAL C 0.4917861D-02

NU= 0.2444940D 01

ITER= 1 PARAMETERS

0.1295068873D-01 0.2183209327D-01

1 NEW RESIDUAL 0.5609383D-04

1 NU IS 0.1222470D 0 1

THE NORM OF THE RELATIVE CORRECTION IS= 0.1370 00

ITER= 2            PARAMETERS  
0.1292835923D-01    0.2199967360D-01  
1 NEW RESIDUAL 0.5468443D-04  
1 NU IS 0.6112350D 00  
THE NORM OF THE RELATIVE CORRECTION IS= 0.663D-02

ITER= 3            PARAMETERS  
0.1287837647D-01    0.2210022751D-01  
1 NEW RESIDUAL 0.5465016D-04  
1 NU IS 0.30561750 00  
THE NORM OF THE RELATIVE CORRECTION IS= 0.439D-02

ITER= 4            PARAMETERS  
0.1286831632D-01    0.2212108054D-01  
1 NEW RESIDUAL 0.5464895D-04  
1 NU IS 0.1528088D 0 0  
THE NORM OF THE RELATIVE CORRECTION IS= 0.905D-03

\*\*\*\*\*  
WEIGHTS

0.37541320 00 0.1936239D 01 -0.1465082D 01

NONLINEAR PARAMETERS

0.1286832D-01 0.2212108D-01