# Sufficiency of parameter invariance conditions in areal and higher order Kawaguchi spaces

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#### Introduction

Let  $V_n$  be an *n*-dimensional space coordinatized by x (Latin indices assume the values 1, ..., n throughout this note). The most elementary variational problem s that of minimizing

$$\int L\{x^i, \dot{x}^i\} dt$$

relative to differentiable curves  $x^i = x^i(t)$  joining given points  $x_0^i$  and  $x_1^i$ . If the integral is to be independent of the curve parametrization, the then Lagrangian L must satisfy the homogeneity equation

(1) 
$$L\{x^i, \lambda \dot{x}^i\} = \lambda L\{x^i, \dot{x}^i\} \text{ for } \lambda > 0.$$

In this note we consider some equations arising from parameter invariance conditions imposed on more general variational problems. The first part is concerned with problems in which the Lagrangian L depends on higher ordered derivatives of the curve  $x^i = x^i(t)$  as in Kawaguchi spaces [1],

(2) 
$$\int L\{x^i, x_1^i, ..., x_m^i\} dt \text{ where } x_r^i = \frac{d^r x^i}{dt^r} \text{ for } r = 1, ..., m.$$

If m=1, then parameter invarience implies the homogeneity condition (1), which in turn implies (we use the summation convention on repeated indices)

$$\frac{\partial L}{\partial x_1^i} x_1^i = L.$$

It was shown by Zermelo [2] and also by Lister [3] in general that the parameter invariance of (2) implies that L satisfies the m partial differential equations (where  $\delta_1^p$  is the Kronecker delta)

(3) 
$$\sum_{r=p}^{m} \frac{r!}{(r-p)!} \frac{\partial L}{\partial x_r^i} x_{r-p+1}^i = \delta_1^p L \quad \text{for} \quad p=1, 2, ..., m.$$

(summation on i=1, ..., n). In this note we show that conversely (3) is sufficient to guarantee the parameter invariance of (2).

The second part of this note deals with the multiple integral problems of the variational calculus, as in Areal spaces [4]. In  $V_n$  an m-dimensional subspace  $S_m$  is defined by the equations  $x^i = x^i(t^\alpha)$ , where Greek indices assume the values 1 to  $m \le n$ . The functions  $x^i(t^\alpha)$  are usually assumed to be of class  $C^2$ , and the multiple integrals defined over  $S_m$  are of the form

(4) 
$$\int_{\mathbb{R}_{+}} \dots \int L\{x^{i}, \dot{x}_{\alpha}^{i}\} dt^{1} \dots dt^{m} \text{ where } \dot{x}_{\alpha}^{i} = \frac{\partial x^{i}}{\partial t^{\alpha}};$$

here the Lagrangian L is a function of the n+nm arguments  $x^i$  and  $\dot{x}^i_{\alpha}$ . Under a transformation  $t^{\alpha}=t^{\alpha}(\tau^{\beta})$  of the Gaussian coordinates of  $S_m$ , the partial derivatives  $\dot{x}^i_{\alpha}$  become  $\lambda^{\beta}_{\alpha}\dot{x}^i_{\beta}$  (summation on  $\beta=1,\ldots,m$ ) where  $\lambda^{\alpha}_{\beta}=\partial t^{\alpha}/\partial \tau^{\beta}$ . If (4) is to be invariant under such transformations then L must satisfy

(5) 
$$L\{x^i, \lambda_\alpha^\beta \dot{x}_\beta^i\} = |\lambda_\alpha^\beta| \cdot L\{x^i, \dot{x}_\alpha^i\}$$

where  $|\lambda_{\alpha}^{\beta}| = \det(\lambda_{\beta}^{\alpha})$  is the Jacobian of the transformation, and should not vanish. It was shown by G. Kobb [5] when n=3, m=2 and generally by W. Grosz [6] pg. 84, that a necessary condition (assuming differentiability etc.) for (5) to hold is that L satisfies the  $m^2$  partial differential equations

$$\frac{\partial L}{\partial \dot{x}^i} \, \dot{x}^i_{\beta} = \delta^{\alpha}_{\beta} L$$

where of course summation on i=1, ..., n is implied. Rund [7] has also shown these equations to be sufficient, but the sufficiency proof required some assumptions on the path of integration (see Martin [4]) and in this paper an alternate sufficiency proof is gigen which avoids these difficulties.

Finally, the equation (5) is a particular case of a matrix functional equation. Let  $[\lambda_{\beta}^{\alpha}]$  denote  $m \times m$  matrices, while  $[x_{\alpha}^{i}]$  denotes rectangular  $m(\text{row}) \times n(\text{column})$  matrices. If f is a real valued function on rectangular  $m \times n$  matrices, then (5) is equivalent to the functional equation

(7) 
$$f\{[\lambda_{\beta}^{\alpha}][x_{\alpha}^{i}]\} = |\lambda_{\beta}^{\alpha}|^{k}f\{[x_{\alpha}^{i}]\}$$

when k=1. If m=1 this is Euler's homogeneity equation. If m=n, so that (7) becomes  $f(XY) = |X|^k f(Y)$  for square matrices X and Y, the equation is readily solved. Let I be the identity matrix; it follows that  $f(X) = f(XI) = |X|^k f(I) = c|X|^k$  for arbitrary c (see Aczél [8] pg. 242). In this note we discuss the general solution of (7) for arbitrary m and n although our results cannot be considered a complete solution of the functional equation.<sup>1</sup>)

## Higher ordered Lagrangians

In this section we show that if L satisfies the m partial differential equations

(3) 
$$\sum_{r=p}^{m} \frac{r!}{(r-p)!} \frac{\partial L}{\partial x_r^i} x_{r-p+1}^i = \delta_1^p L \quad \text{for} \quad p=1, ..., m$$

<sup>1)</sup> As pointed out by Prof. Aczél, equation (7) characterizes the scalar density concomitans of weight k defined on n vectors.

where L is a function of the n(m+1) variables,  $L\{x^i, x^i_1, ..., x^i_m\}$ , then

$$L\left\{x^{i}, \frac{dx^{i}}{dt}, \dots, \frac{d^{m}x^{i}}{dt^{m}}\right\} = \frac{d\tau}{dt}L\left\{x^{i}, \frac{dx^{i}}{d\tau}, \dots, \frac{d^{m}x^{i}}{d\tau^{m}}\right\}$$

for any parameter transformation  $\tau = \varphi(t)$  of class  $C^m$  satisfying  $\varphi'(t) > 0$ ; we assume that L has continuous partial derivatives for all  $x_r^i$ , except perhaps for  $(x_1^1)^2 + ... + ... + (x_1^n)^2 = 0$ . (This latter restriction is imposed for simplicity and can be considerably weakened, as indicated below, if so desired).

The following combinatorial formulas will be needed. Let  $A_{\sigma}^{k}(m)$  denote

(8.1) 
$$A_{\sigma}^{k}(m) = \begin{cases} 0 & \text{if} \\ \sum_{(q)} \frac{1}{q_{1}! \dots q_{m}!} \lambda_{1}^{q_{1}} \dots \lambda_{m}^{q_{m}} & \text{if} \quad \sigma \geq k \end{cases}$$

where the summation is over integral  $q_i \ge 0$  satisfying

(8.2) 
$$q_1 + q_2 + \dots + q_m = k$$
$$q_1 + 2q_2 + \dots + mq_m = \sigma.$$

The  $A_{\sigma}^{k}(m)$  are essentially multinomial coefficients [9], satisfying

(9) 
$$(\lambda_1 z + \lambda_2 z^2 + \dots + \lambda_m z^m)^k = k! \sum_{\sigma} A_{\sigma}^k(m) z^{\sigma}.$$

From (9) it follows (by differentiating with respect to  $\lambda_p$ ) that

(10) 
$$\frac{\partial A_{\sigma}^{k}(m)}{\partial \lambda_{p}} = \begin{cases} 0 & \text{if } \sigma < p+k-1, \\ A_{\sigma-p}^{k-1}(m) & \text{if } \sigma \ge k-p+1; \end{cases}$$

and also (differentiating with respect to z)

(11) 
$$\sum_{\alpha=1}^{\sigma-k+2} \alpha \lambda_{\alpha} A_{\sigma-\alpha+1}^{k-1}(m) = (\sigma+1) A_{\sigma+1}^{k}(m).$$

Since (3) is to hold for arbitrary variables  $x_r^i$ , then also for  $u_r^i$  related to the  $x_r^i$  by

(12) 
$$u_r^i = \sum_{s=1}^r x_s^i r! A_r^s(m).$$

Replacing  $x_r^i$  throughout (3) by  $u_r^i$ , the system of equations becomes

(13) 
$$\Omega_{p} \stackrel{\text{def}}{=} \sum_{r=p}^{m} \frac{\partial L}{\partial u_{r}^{i}} r! \sum_{s=1}^{r-p+1} x_{s}^{i} (r-p+1) A_{r-p+1}^{s} (m) = \delta_{1}^{p} L(u_{r}^{i}),$$

where the left side defines the m quantities  $\Omega_p$ .

If  $\lambda_1 \neq 0$ , we wish to show that by forming linear combinations of the equations (13), one can obtain the system

$$\frac{\partial L(u_r^i)}{\partial \lambda_p} = \frac{1}{\lambda_1} L(u_r^i) \frac{\partial \lambda_1}{\partial \lambda_p} \qquad p = 1, ..., m$$

or, expanding and using (10) and (12),

However the required linear combinations of the equations  $\Omega_p = \delta_1^p L$  which yield  $\Lambda_p = \frac{1}{\lambda_1} \delta_1^p L$  are relatively complicated and it is easier to reverse the process as follows.

Let  $X_p$  and  $Y_p$  be arbitrary variables, p = 1, ..., m, related by

(16) 
$$X_{p} = \sum_{\sigma=1}^{m-p+1} \sigma \lambda_{\sigma} Y_{p+\sigma-1} \qquad p=1, ..., m.$$

Then it is easily verified that, given  $\lambda_1 \neq 0$ , the  $Y_p$  are given by

(17) 
$$Y_{p} = \sum_{\sigma=1}^{m-p+1} Q_{p}^{\sigma}(\lambda_{r}) X_{p+\sigma-1} \qquad p=1, ..., m$$

where the  $Q_p^{\sigma}(\lambda_r)$  are polynomials in the  $\lambda_r$  except for the presence of  $\lambda_1$  to negative powers. In particular

$$Y_m = \frac{1}{\lambda_1} X_m; \quad Y_{m-1} = \frac{1}{\lambda_1} X_{m-1} - 2 \frac{\lambda_2}{\lambda_1^2} X_m; \quad \text{etc.}$$

Instead of applying (17) to equations (13) and obtaining (15), it is sufficient to apply (16) to equations (15) and verifying that (13) results. Explicitly, we will show that

(18.1) 
$$\Omega_p = \sum_{\sigma=1}^{m-p+1} \sigma \lambda_{\sigma} \Lambda_{p+\sigma-1} \qquad p=1, ..., m$$

(18.2) 
$$\delta_1^p L = \sum_{\sigma=1}^{m-p+1} \sigma \lambda_{\sigma} \left( \frac{1}{\lambda_1} \delta_1^{p+\sigma-1} L \right) \qquad p = 1, ..., m.$$

It will then follow that equations (15) are a consequence of (13) whenever  $\lambda_1 \neq 0$ . Equation (18. 2) is readily verified since the right hand side is zero unless  $p+\sigma=2$ ; but since  $p \geq 1$  and  $\sigma \geq 1$ , this can only happen for  $p=\sigma=1$ , as required. To verify (18. 1), substituting from the definition of  $\Lambda_p$  given in (15), one obtains

$$\sum_{\sigma=1}^{m-p+1} \sum_{r=p+\sigma-1}^{m} \sum_{s=1}^{r-(p+\sigma-1)+1} \frac{\partial L}{\partial u^i_r} \, x^i_s \, r! \, \sigma \lambda_\sigma \, A^{s-1}_{r-(p+\sigma-1)},$$

and interchanging orders of summation yields

$$\sum_{r=p}^{m} \frac{\partial L}{\partial u_r^i} r! \sum_{s=1}^{r-p+1} x_s^i \sum_{\sigma=1}^{(r-p)-s+2} \sigma \lambda_{\sigma} A_{(r-p)-\sigma+1}^{s-1}.$$

By (11), this is  $\Omega_p$  as required.

It follows that applying (17) to (13) will yield (14) whenever  $\lambda_1 \neq 0$ . But (14) can be integrated since we may write it in the form

$$\left\{\lambda_1 \frac{\partial L}{\partial \lambda_p} - \frac{\partial \lambda_1}{\partial \lambda_p} L\right\} \lambda_1^{-2} = \frac{\partial}{\partial \lambda_p} \left\{\lambda_1^{-1} L(u_r^i)\right\} = 0.$$

An integration along a curve<sup>2</sup>) in the  $\lambda_r$  space from  $\lambda_r = \delta_r^1$  to an arbitrary  $\lambda_r$ , avoiding  $\lambda_1 = 0$ , yields

(19) 
$$L(u_r^i) = \lambda_1(x_r^i) \quad \text{for} \quad \lambda_1 \neq 0$$

since, by (8.1) and (12),  $u_r^i$  reduces to  $x_r^i$  when  $\lambda_r = \delta_r^1$ . But (19) is the required formula for parameter invariance since, by setting  $\lambda_r = \frac{1}{r!} \frac{d^r \tau}{dt^r}$ , equation (19), when fully expanded, becomes

$$L\left\{\sum_{s=1}^{r} \frac{d^{s} x}{d\tau^{s}} r! \sum_{(q)} \frac{1}{q_{1}! \dots q_{m}!} \left(\frac{\frac{d\tau}{dt}}{1!}\right)^{q_{1}} \dots \left(\frac{\frac{d^{m} \tau}{dt^{m}}}{m!}\right)^{q_{m}}\right\} = \frac{d\tau}{dt} L\left\{\frac{d^{r} x^{i}}{d\tau^{r}}\right\}.$$

But the arguments on the left are Faa di Bruno's formula [10] for the rth derivative of the composite function  $x\{\Phi(t)\}$ , verifying the parameter invariance q. e. d.

## Multiple integrals

In this case the Lagrangian  $L\{x^i, \dot{x}^i_\alpha\}$  is a function of the n+nm variables  $\dot{x}^i_{\alpha} = \partial x^i / \partial t^{\alpha}$  and  $x^i$  where i = 1, ..., n and  $\alpha = 1, ..., m \le n$ . As indicated in the introduction we wish to show that if  $L\{x^i, \dot{x}^i_{\alpha}\}$  satisfies the  $m^2$  partial differential equations

(6) 
$$\frac{\partial L}{\partial \dot{x}^i} \dot{x}^i_{\beta} = \delta^{\alpha}_{\beta} L$$

then L satisfies the functional equation

(5) 
$$L\{x^i, \lambda_\alpha^\beta \dot{x}_\beta^i\} = |\lambda_\beta^\alpha| L\{x^i, \dot{x}_\alpha^i\} \text{ for } |\lambda_\beta^\alpha| \neq 0$$

and hence leads to parameter invarient multiple integrals.

This has already been shown by RUND [7], and we briefly outline his method in order to compare it with the method of this note. Let  $[\lambda_{\beta}^{\alpha}]$  be an arbitrary  $m \times m$ matrix and set  $J = |\lambda_{\beta}^{\alpha}|$ . Assume  $J \neq 0$  and let  $\Lambda_{\alpha}^{\beta}$  denote the cofactor of  $\lambda_{\beta}^{\alpha}$  in J. Replace  $\dot{x}_{\beta}^{i}$  throughout (6) by  $u_{\beta}^{i} = \lambda_{\beta}^{\alpha} \dot{x}_{\gamma}^{i}$  and multiply through by  $\Lambda_{\mu}^{\beta}$ . The result may be written in the form

$$J\frac{\partial L}{\partial u_{\nu}^{i}}\dot{x}_{\nu}^{i}=\Lambda_{\nu}^{\alpha}L.$$

<sup>2)</sup> We are using the fact that  $\int_{0}^{P_2} \frac{\partial F}{\partial \lambda_r} d\lambda_r = F|_{P_1}^{P_2}$ , (summation on r). Here  $P_2$  is the point with

coordinates  $[\lambda_1 \neq 0, \lambda_2, \dots, \lambda_m]$  and  $P_1$  has coordinates  $[1, 0, \dots, 0]$ . We are assuming that the points  $P_1$  and  $P_2$  can be joined by a curve along which  $\lambda_1 \neq 0$ , and for which F is continuously differentiable. (This is the condition which should replace the original condition on L concerning differentiability etc.).

But  $\partial J/\partial \lambda_{\alpha}^{\nu} = \Lambda_{\nu}^{\alpha}$  while  $\partial L/\partial \lambda_{\alpha}^{\nu} = (\partial L/\partial u_{\alpha}^{i})x_{\nu}^{i}$  so that this becomes

$$\frac{\partial}{\partial \lambda_{\alpha}^{\nu}} \{J^{-1} L\} = 0$$
 provided  $J \neq 0$ .

Rund now integrates this equation along a path in the  $\lambda_{\beta}^{\alpha}$  space from the initial point  $\lambda_{\beta}^{\alpha} = \delta_{\beta}^{\alpha}$  to an arbitrary  $\lambda_{\beta}^{\alpha}$ , assuming that along the path  $J \neq 0$ . Since J = 1 at the initial point, the integrated equation becomes (5).

Martin [4], pg. 120, has raised some questions concerning this procedure, mainly concerning the path of integration: is every point in  $\lambda_{\beta}^{\alpha}$  space, even if  $J = |\lambda_{\beta}^{\alpha}| > 0$ , connected to the point  $\delta_{\beta}^{\alpha}$  by a path along which  $J \neq 0$ ? We will not attempt to

answer this questions here.

In this note, the equation (5) is derived from (6) in a way which avoids such problems in the  $\lambda_{\beta}^{\alpha}$  space. Instead we assume simply that (6) holds in a region of the  $\dot{x}_{\beta}^{i}$  space for which the rank of the matrix  $[\dot{x}_{\beta}^{i}]$  is m. Since the system (6) is not affected by interchanging the columns of the matrix  $[\dot{x}_{\beta}^{i}]$  we assume for simplicity that the variables have been re-labeled such that the first m columns  $[\dot{x}_{\beta}^{i}]$  have non-vanishing determinant in the region in question. We write the matrix in the partitioned form

$$[\dot{x}^i_{\alpha}] = [\dot{x}^{\alpha}_{\beta} : \dot{x}^p_{\beta}]$$

where p = m+1, ..., n. (Exception is made if m=n. This case is simple and will be treated below). Since  $\dot{x}_{\beta}^{\alpha} \neq 0$  by assumption, let  $[X_{\beta}^{\alpha}] = [\dot{x}_{\beta}^{\alpha}]^{-1}$ , the inverse matrix. We now perform the following change of variables on the system (6):

(20) 
$$\begin{cases} \text{leave } \dot{x}_{\beta}^{\alpha} \text{ unchanged} \\ \text{replace } \dot{x}_{\beta}^{p} \text{ by } u_{\beta}^{p} = X_{\beta}^{\gamma} \dot{x}_{\gamma}^{p}, \text{ for } p = m+1, ..., n. \end{cases}$$

Since the matrix  $[X_{\beta}^{\gamma}]$  has an inverse, this transformation is one-one. Finally, replace  $L(x^i, \dot{x}_{\alpha}^i)$  by

(21) 
$$L(x^{i}, \dot{x}_{\alpha}^{i}) = |\dot{x}_{\beta}^{\alpha}| l\{x^{i}, \dot{x}_{\beta}^{\alpha}, u_{\beta}^{p}\}.$$

Since l is arbitrary and  $|\dot{x}^{\alpha}_{\beta}| \neq 0$ , this places no restriction on the system. In the remainder of this section it will be shown that under the transformation (20) and (21) the system of equations (6) are equivalent to the system

$$\frac{\partial l}{\partial \dot{x}_{\alpha}^{\nu}} = 0$$

so that the general solution of (6) is given by

(23) 
$$L\{x^{i}, \dot{x}_{\alpha}^{i}\} = |\dot{x}_{\beta}^{\alpha}| I\{x^{i}, X_{\beta}^{\gamma} \dot{x}_{\gamma}^{p}\} \qquad p = m+1, ..., n,$$

for arbitrary l. In the case m = n, this reduces to

$$L\{x^i, \dot{x}^i_j\} = |\dot{x}^i_j| \, l\{x^i\}.$$

Further, any function L of the form (23) satisfies the functional equation (5), given  $|\lambda_B^{\alpha}| \neq 0$ , as was required.

In order to verify the assertion, let  $\Delta = |\dot{x}_{\beta}^{\alpha}|$ . Then, if  $C_{\alpha}^{\beta}$  denotes the cofactor of  $\dot{x}_{\beta}^{\alpha}$  in  $\Delta$ , it is clear that

(24) 
$$\frac{\partial \Delta}{\partial \dot{x}_{\gamma}^{\alpha}} \dot{x}_{\beta}^{\gamma} = C_{\gamma}^{\alpha} \dot{x}_{\beta}^{\gamma} = \delta_{\beta}^{\alpha} \Delta.$$

Also, since  $X^{\sigma}_{\mu}\dot{x}^{\mu}_{\tau} = \delta^{\sigma}_{\tau}$ , it follows that

$$\frac{\partial X^{\sigma}_{\mu}}{\partial \dot{x}^{\gamma}_{\alpha}}\dot{x}^{\mu}_{\tau} = -X^{\sigma}_{\mu}\frac{\partial \dot{x}^{\mu}_{\tau}}{\partial \dot{x}^{\gamma}_{\alpha}} = -X^{\sigma}_{\gamma}\,\delta^{\alpha}_{\tau},$$

and multiplying by  $X_{\nu}^{\tau} \hat{x}_{\beta}^{\gamma}$  yields

(25) 
$$\frac{\partial X_{\nu}^{\sigma}}{\partial \dot{x}_{\alpha}^{\gamma}} \dot{x}_{\beta}^{\gamma} = -X_{\gamma}^{\sigma} X_{\nu}^{\alpha} \dot{x}_{\beta}^{\gamma} = -\delta_{\beta}^{\sigma} X_{\nu}^{\alpha}.$$

We now write (6) in the form

$$rac{\partial L}{\partial \dot{x}^{\gamma}_{lpha}}\dot{x}^{\gamma}_{eta}+rac{\partial L}{\partial \dot{x}^{p}_{lpha}}\dot{x}^{p}_{eta}=\delta^{lpha}_{eta}\,L$$

summation on  $\gamma = 1, ..., m$  and p = m + 1, ..., n. Replacing L by  $L = \Delta l$  as in (21), this equation becomes

(26) 
$$\frac{\partial \Delta}{\partial \dot{x}_{\alpha}^{\gamma}} \dot{x}_{\beta}^{\gamma} l + \Delta \frac{\partial l}{\partial \dot{x}_{\alpha}^{\gamma}} \dot{x}_{\beta}^{\gamma} + \Delta \frac{\partial l}{\partial u_{\nu}^{p}} \frac{\partial X_{\nu}^{\sigma}}{\partial \dot{x}_{\alpha}^{\gamma}} \dot{x}_{\beta}^{\gamma} \dot{x}_{\sigma}^{p} + \Delta \frac{\partial l}{\partial u_{\nu}^{p}} X_{\nu}^{\alpha} \dot{x}_{\beta}^{p} = \delta_{\beta}^{\alpha} \Delta l.$$

Using (24), the first term cancels with the right hand side, while by (25) the third and fourth terms cancell. Since  $\Delta \neq 0$ , this equation reduces to

$$\frac{\partial l}{\partial \dot{x}_{\beta}^{\gamma}} \dot{x}_{\beta}^{\gamma} = 0$$

and multiplying by  $X_{\nu}^{\beta}$  yields (22). If m=n, then the variables  $u_{\beta}^{p}$  are missing in (26), but the result is the same. It remains to show that if L is of the form specified in (23) then the functional equation (5) is satisfied. This will be shown with more generality in the next section.

## The matrix functional equation

In this section we discuss a generalization of the functional equation treated (with differentiability assumptions) in the previous section. As before  $[\lambda_{\beta}^{\alpha}]$  denotes an  $m \times m$  matrix,  $[x_{\alpha}^{i}]$  an  $m(\text{row}) \times n(\text{column})$  matrix where now n may be less than m, and f is a real valued function on rectangular  $m \times n$  matrices, satisfying

(7) 
$$f\{[\lambda_{\beta}^{\alpha}][x_{\alpha}^{i}]\} = |\lambda_{\beta}^{\alpha}|^{k}f\{[x_{\alpha}^{i}]\}$$

We do not give a complete solution of (7) since this seems to require many distinctions in the possible ways in which the rank of the matrix  $[x_{\alpha}^{i}]$  can be m. This is

analogous to the case  $F\{\lambda x, \lambda y\} = \lambda^k F(x, y)$  which requires the four distinct cases (ACZÉL [8] p. 160):  $y \neq 0$ , y = 0 and  $x \neq 0$ , y = x = 0 and  $k \neq 0$ , y = x = 0 and k = 0. Here, the matrix  $[x_\alpha^i]$  can have rank m in  $\frac{n!}{m!(n-m)!}$  ways (the number of ways of choosing m distinct columns from the total of n columns).

In this section we merely show that

(27) 
$$f\{[x_{\alpha}^{i}]\} = |x_{\beta}^{\alpha}|^{k} F\{X_{\beta}^{\gamma} x_{\gamma}^{p}\}, \quad |x_{\beta}^{\alpha}| \neq 0$$

for arbitrary F if the rank of the matrix  $[x_{\beta}^{\alpha}:x_{\beta}^{p}]$  is m and  $|x_{\beta}^{\alpha}|$  is the non-vanishing determinant, while

$$f\{[x_{\alpha}^i]\} = 0$$

if the rank of the matrix  $[x_{\alpha}^{i}]$  is less than m, while  $k \neq 0$ . Here  $X_{\beta}^{\gamma}$  and  $x_{\gamma}^{p}$  have the same interpretation<sup>3</sup>) as in the previous section, and F is an arbitrary function of the (n-m)m variables  $X_{\beta}^{\gamma}x_{\gamma}^{p}$ . In the case m=1 equation (27) reads.

$$f\{x^1, x^2, ..., x^n\} = (x^1)^k F\left\{\frac{x^2}{x^1}, ..., \frac{x^n}{x^1}\right\}, \qquad x^1 \neq 0.$$

The above assertion may be verified as follows. If the rank of  $[x_{\alpha}^i]$  is less than m, then  $[x_{\alpha}^i]$  is row-equivalent to a matrix with zero first row. Alternatively there exists a non singular matrix  $[\lambda_{\beta}^{\alpha}]$  such that the first row in  $[\lambda_{\beta}^{\alpha}][x_{\alpha}^i]$  is zero. But then if  $[\sigma_{\beta}^{\alpha}]$  is the matrix

$$[\sigma_{\beta}^{\alpha}] = \begin{bmatrix} \sigma & 0 \dots 0 \\ 0 & 1 \dots 0 \\ \vdots & \vdots & \vdots \\ 0 & 0 \dots 1 \end{bmatrix},$$

the matrix  $[\sigma_{\beta}^{\alpha}][\lambda_{\beta}^{\alpha}][x_{\alpha}^{i}]$  is independent of  $\sigma$ . It follows by substituting in (7) that the right hand side of the equation

$$f\{[\sigma_{\beta}^{\alpha}][\lambda_{\beta}^{\alpha}][x_{\alpha}^{i}]\} = \sigma^{k}|\lambda_{\beta}^{\alpha}|^{k}f\{[x_{\alpha}^{i}]\}, \qquad k \neq 0$$

is independent of  $\sigma$  and since  $|\lambda_{\beta}^{\alpha}| \neq 0$ , this implies  $f\{[x_{\alpha}^{i}]\} = 0$ .

If the matrix  $[x_{\beta}^{i}] = [x_{\beta}^{\alpha} : x_{\beta}^{p}]$  has rank m (which implies that  $n \ge m$ ), and if further  $|x_{\beta}^{\alpha}| \ne 0$ , the  $[\lambda_{\beta}^{\alpha}]$  in (7) can be chosen as the inverse matrix  $[X_{\beta}^{\alpha}] = [x_{\beta}^{\alpha}]^{-1}$ , resulting in the equation

$$f\{[X_{\beta}^{\alpha}][x_{\beta}^{\alpha}:x_{\beta}^{p}]\} = f\{[\delta_{\beta}^{\alpha}:X_{\beta}^{\gamma}x_{\gamma}^{p}]\} = |X_{\beta}^{\alpha}|^{k}f\{[x_{\alpha}^{i}]\}.$$

Since  $|X_{\beta}^{\alpha}|^k = |x_{\beta}^{\alpha}|^{-k}$  it follows that f has the form specified in (27). Conversely, if f is given by (27) then

(28) 
$$f\{[\lambda_{\beta}^{\alpha}][x_{\alpha}^{i}]\} = f\{[\lambda_{\beta}^{\gamma} x_{\gamma}^{\alpha} : \lambda_{\beta}^{\gamma} x_{\gamma}^{p}]\} = |\lambda_{\beta}^{\gamma} x_{\gamma}^{\alpha}|^{k} F\{Y_{\beta}^{\gamma} y_{\gamma}^{p}\}$$

<sup>3)</sup> This is but a slight generalization of the results of W. Grosz [6], pg. 82—83, for k=1. In his notation  $p_0 = |x_{\beta}^{\alpha}|$  and  $p_{h,k} = -p_0 X_h^{\gamma} x_{\gamma}^{m+k}$ .

where 
$$y_{\gamma}^{p} = \lambda_{\gamma}^{v} x_{v}^{p}$$
 and  $[Y_{\beta}^{\gamma}] = [\lambda_{\beta}^{v} x_{v}^{\gamma}]^{-1} = [x_{\beta}^{\alpha}]^{-1} [\lambda_{\beta}^{\alpha}]^{-1}$ . Therefore  $[Y_{\beta}^{\gamma} y_{\gamma}^{p}] = [X_{\beta}^{\alpha}] [\lambda_{\beta}^{\alpha}]^{-1} [\lambda_{\beta}^{\alpha}] [x_{\beta}^{p}] = [X_{\beta}^{\gamma} x_{v}^{p}]$ .

Hence (28) becomes

$$f\{[\lambda_{\beta}^{\alpha}][x_{\alpha}^{i}]\} = |\lambda_{\beta}^{\alpha}|^{k} |x_{\beta}^{\alpha}|^{k} F\{X_{\beta}^{\gamma} x_{\gamma}^{p}\} = |\lambda_{\beta}^{\alpha}|^{k} f\{[x_{\alpha}^{i}]\}$$

which verifies the assertion that the function f in (27) satisfies (7), and also completes the proof of the previous section.

### References

- [1] A complete bibliography on Kawaguchi spaces up to 1950, compiled by H. Schubert, appears in the appendix to the reprint of Finsler's thesis, "Über Kurven und Flächen in allgemeinen Raumen" (1918), unveränderter Nachdruck, Basel, 1951.
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